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## CORSO DI DOTTORATO DI RICERCA IN: INGEGNERIA MECCATRONICA E DELL'INNOVAZIONE MECCANICA DEL PRODOTTO <br> CURRICOLO: MECCATRONICA <br> CICLO XXIX <br> EIGENSTRUCTURE ASSIGNMENT IN VIBRATING SYSTEMS <br> THROUGH ACTIVE AND PASSIVE APPROACHES

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## Abstract

The dynamic behaviour of a vibrating system depends on its eigenstructure, which consists of the eigenvalues and the eigenvectors. In fact, eigenvalues define natural frequencies, damping and settling time, while eigenvectors define the spatial distribution of vibrations, i.e. the mode shape, and also affect the sensitivity of eigenvalues with respect to the system parameters.

Therefore, eigenstructure assignment, which is aimed at modifying the system in such a way that it features the desired set of eigenvalues and eigenvectors, is of fundamental importance in mechanical design. However, similarly to several other inverse problems, eigenstructure assignment is inherently challenging, due to its ill-posed nature. Despite the recent advancements of the state of the art in eigenstructure assignment, in fact, there are still important open issues.

The available methods for eigenstructure assignment can be grouped into two classes: passive approaches, which consist in modifying the physical parameters of the system, and active approaches, which consist in employing actuators and sensors to exert suitable control forces as determined by a specified control law. Since both these approaches have advantages and drawbacks, it is important to choose the most appropriate strategy for the application of interest.

In the present thesis, in fact, are collected passive, active, and even hybrid methods, in which active and passive techniques are concurrently employed. All the methods proposed in the thesis are aimed at solving open issues that emerged from the literature and which have applicative relevance, as well as theoretical. In contrast to several state-of-the-art methods, in fact, the proposed ones implement strategies that enable to ensure that the computed solutions are meaningful and feasible. Moreover, given that in modern mechanical design large-scale systems are increasingly common, computational issues have become a major concern and thus have been adequately addressed in the thesis.

The proposed methods have been developed to be general and broadly applicable. In order to demonstrate the versatility of the methods, in the thesis it is provided an extensive numerical assessment, hence diverse test-cases have been used for validation purposes. In order to evaluate without bias the performances of the proposed methods, it has been chosen to employ well-
established benchmarks from the literature. Moreover, selected experimental applications are presented in the thesis, in order to determine the capabilities of the developed methods when critically challenged.

Given the focus on these issues, it is expected that the methods here proposed can constitute effective tools to improve the dynamic behaviour of vibrating systems and it is hoped that the present work could contribute to spread the use of eigenstructure assignment in the solution of engineering design problems.

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## Chapter 1. Introduction

### 1.1. Scope of the thesis

Vibration occurs in every mechanical system that have elastic forces, and hence a proper understanding and mastery of such a phenomenon is of primary importance in mechanical design. Vibrations might be unwanted, being a major cause of stress and fatigue and reducing precision and accuracy of machines, or desired, since required by the processes of interest. In both cases, it is often necessary to comply with design specifications that concerns the vibratory behaviour of mechanical systems.

The most established model for the analysis of multiple degrees of freedom vibrating systems is the well-known second order linear model

$$
\begin{equation*}
\mathbf{M} \ddot{\boldsymbol{q}}(t)+\mathbf{C} \dot{\boldsymbol{q}}(t)+\mathbf{K} \boldsymbol{q}(t)=\boldsymbol{f}(t) \tag{1.1}
\end{equation*}
$$

where $\mathbf{M}, \mathbf{C}, \mathbf{K} \in \mathbb{R}^{N \times N}$ are the mass, damping and stiffness matrices, respectively, $\boldsymbol{q} \in \mathbb{R}^{N}$ is the displacements vector and $\boldsymbol{f} \in \mathbb{R}^{N}$ represents the external forces that are exerted on the system (e.g. disturbances, control forces). From the theory of linear ordinary differential equations, it follows that the dynamical behaviour of a vibrating system is determined by the eigenvalues $\lambda \in \mathbb{C}$ and the eigenvectors $\boldsymbol{u} \in \mathbb{C}^{N}$ associated to the system, namely the solutions of the quadratic eigenproblem

$$
\begin{equation*}
\left[\lambda^{2} \mathbf{M}+\lambda \mathbf{C}+\mathbf{K}\right] \boldsymbol{u}=\mathbf{0} \tag{1.2}
\end{equation*}
$$

The pair $(\lambda, \boldsymbol{u})$ is called eigenpair and it can be easily proved that the set of eigenpairs consists of $2 N$ complex self-conjugate elements. Eigenvalues define the natural frequency (or eigenfrequency) $\omega$ and the damping factor $\zeta$ of each mode of vibration, being such quantities related by

$$
\begin{equation*}
\lambda=-\zeta \omega \pm i \omega \sqrt{1-\zeta^{2}} \tag{1.3}
\end{equation*}
$$

Eigenvectors instead define the mode shapes, namely the spatial distribution of vibrations over the system. Moreover, eigenvectors affect the sensitivity of eigenvalues with respect to system
parameters, and hence robustness. In fact, let us consider a parameter $x$, then the following formula [1] expresses the sensitivity of an eigenvalue $\lambda$ with respect to $x$ :

$$
\begin{equation*}
\frac{\partial \lambda}{\partial x}=\lambda \frac{\boldsymbol{u}^{T}\left[\frac{\partial \mathbf{K}}{\partial x}-\lambda^{2} \frac{\partial \mathbf{M}}{\partial x}+i \lambda \frac{\partial \mathbf{C}}{\partial x}\right] \boldsymbol{u}}{\boldsymbol{u}^{T}\left[\lambda^{2} \mathbf{M}+\mathbf{K}\right] \boldsymbol{u}} \tag{1.4}
\end{equation*}
$$

where $\boldsymbol{u}$ is the eigenvector associate to $\lambda$.

At this point it is worth to mention the special case in which the system is undamped, i.e. $\mathbf{C}=\mathbf{0}$, the mass matrix is symmetric and positive definite and the stiffness matrix is symmetric and positive semi-definite. Such assumptions are reasonable, since they allow to accurately represent most of the vibrating system of interest. Moreover, the mathematical complexity of the problem can be significantly simplified. In fact, in this case it holds

$$
\begin{equation*}
\lambda= \pm i \omega \tag{1.5}
\end{equation*}
$$

which is exactly (1.3) for $\zeta=0$. Hence it is convenient to solve the equivalent linear eigenproblem

$$
\begin{equation*}
\left[\omega^{2} \mathbf{M}-\mathbf{K}\right] \boldsymbol{u}=\mathbf{0} \tag{1.6}
\end{equation*}
$$

In such a case, there are $N$ real and non-negative eigenvalues, each one associated to a real eigenvector.

Regardless of the employed system model and the resulting eigenproblem to be solved, it is evident that the eigenvalues and eigenvectors, which together constitute the eigenstructure, strongly characterize the dynamic behaviour of the system. The design of vibrating systems aimed at satisfying specifications on eigenvalues and eigenvectors, which is commonly known as eigenstructure assignment (EA), has drawn increasing interest over the last few decades. The most natural mathematical framework for such problems is constituted by the inverse eigenproblems, which consist in the determination of the system model that features a desired set of eigenvalues and eigenvectors, indeed. The development of effective and reliable numerical methods is inherently challenging, due to the ill-posed nature of this kind of inverse problems. In fact, despite the remarkable advancements made over the recent years, the state of the art is still unable to solve some relevant design problems.

In the present thesis are collected novel techniques for eigenstructure assignment, whose aim is to solve open problems emerged from the literature that reflect actual design needs. It has been chosen to develop methods as general as possible, thus allowing a broad scope of application. Feasibility of the computed solutions, although neglected by several methods available in the literature, has been considered a mandatory requirement since it is often necessary to comply with technical and economic constraints. Moreover, attention has been paid to computational issues, such as numerical conditioning and CPU time, which are critical especially in case of large-scale models. In fact, for an accurate approximation of complex mechanical systems, typically by means of Finite Element models, it is required to increase substantially the number of degrees of freedom of the model and, despite the considerable capabilities of modern computers, computational issues can still be concerning. Finally, an extensive numerical validation of the proposed methods has been carried out. The chosen benchmarks are taken from the literature, allowing an unbiased assessment of the performances achieved.

### 1.2. Literature review

The approaches to eigenstructure assignment can be basically divided into passive control and active control.

Passive control, also known as Dynamic Structural Modification (DSM), aims at achieving the desired eigenstructure by modifying the physical properties of the system (typically the inertial and elastic parameters). From the mathematical point of view, the matrices in (1.1) are altered, usually by additive modification, in such a way that the system features the desired eigenstructure. Such an approach is often chosen because it does not require neither electronic devices nor external power. Moreover, the resulting system is assured to be stable, in fact the system modifications are inherently symmetric. However, in practice, not every desired specification can be obtained, because there are constraints on the modifications that naturally arise from feasibility considerations. Therefore, passive control is not always capable of assigning the desired eigenstructure with sufficient accuracy.

Active control consists in employing a feedback closed-loop scheme in such a way that the controlled system features the desired eigenstructure. In practice, suitable external forces are
exerted on the system by controlled actuators. Although such an approach requires the use of electronic devices, active control is preferable when structural modifications are impractical or inconvenient. Moreover, active control is capable of stabilizing unstable systems and it ensures adequate performances regardless of the operating conditions, even in the presence of external disturbances. Active vibration control relies on a solid theoretical background. For the assignment of just the eigenvalues, i.e. the so-called "pole placement", several methods based on state-feedback are available. The assignment of eigenvectors, in contrast, is affected by an inherent limitation that restricts the range of eigenvectors that can be actually obtained.

Both passive and active control have advantages and shortcomings. Recently in the literature appeared also hybrid methods that combine these two approaches, in order to achieve superior performances with respect to the sole employment of either passive or active control, but the state of the art is still undeveloped. In this Section, a brief literature review of passive and active control is given.

### 1.2.1. Passive control

Passive control, or Dynamic Structural Modification, is related to two distinct problems. The direct problem is the prediction of the modified natural frequencies and mode shapes when the employed perturbation of the system is known. In contrast, the inverse problem is the computation of the appropriate values of the physical parameters of the system that results in the desired dynamic behaviour. In such a sense, inverse structural modification is a suitable technique for eigenstructure assignment, as demonstrated by numerous paper appeared in literature. Such methods can be used either for the design of new vibrating systems or for the optimization of existing ones. In both the cases, some knowledge of the original system model, although approximated, is required. The criterion that is chosen to classify the methods for passive control is the employed description of the system model. Hence the works in the literature can be organized as follows:

- the approaches based on the physical model of the system, i.e. its mass, damping and stiffness matrices;
- the approaches based on modal data, i.e. a (not necessarily complete) set of eigenpairs;
- the approaches based on the Frequency Response Functions, and in particular the system receptances.


## Approaches based on physical models

Given the popularity of Finite Element modelling in modern mechanical design, it is not surprising that several methods for passive control rely on the physical system model. This approach is suitable at the design stage of new mechanical systems, in which are often employed Finite Element models. The earliest techniques arise from the investigation of eigenvalue and eigenvector sensitivity to system parameters. For example, in [2], Farahani and Bahai approximate the modification of the system eigenstructure through its first derivative with respect to the parameter modification. In [3] the same authors extend the previous method employing the second order Taylor expansion. Such an approach is effective and computationally efficient as long as the desired changes are small, otherwise it is necessary to perform multiple iterations, affecting the reliability and the convergence of the method.

Liangsheng in [4] proposes to simultaneously modify the mass and stiffness matrices for assigning a subset of the modes. Arbitrary topologies of the matrices are allowed, and the structural modifications are conveniently computed solving a linear system. The number of assigned modes, however, is not arbitrary as it depends on the model dimension and the number of modifiable parameters. Moreover, no constraints are taken into account to ensure feasibility of the modifications.

Wang et al., in [5], manage to maximize the fundamental frequency of beams and plates by a proper positioning of elastic or rigid supports. An iterative procedure, which relies on the sensitivity of the frequency with respect to the support position, is developed and a strategy to improve convergence is proposed. The minimum stiffness of an elastic support that enables to maximize the fundamental frequency is investigated by Wang et al. in [6], where analytical solutions for different boundary conditions are given.

In [7] Ram determines the distribution of mass that enables to maximize the lowest eigenvalue or to minimize the highest. Closed-form solutions are provided for an axially vibrating rod and an Euler-Bernoulli beam. The problem is cast as a quadratic program with quadratic constraints (the latter to constrain the total mass) and a recursive strategy of solution is proposed.

The work of Richiedei et al. [8] aims to handle very general eigenstructure assignment tasks, allowing to assign an arbitrary number of modes and dealing with matrices with arbitrary topology. The method exploits convex optimization to ensure that the computed solution is globally optimal. Moreover, constraints on the modification of the physical parameters can be considered. Such an approach has been later extended by Ouyang et al. [9] that recast the problem as a mixed-integer optimization problem to take into account discrete modifications, in order to reflect the industrial need of modularisation.

## Approaches based on modal data

Although the mass, damping and stiffness matrices represent a simple and powerful way to model vibrating systems, the measurement of the physical quantities that appear in such matrices is often impractical and prone to errors. Nevertheless, the sole eigenstructure of the original system, without the underlying physical matrices, has been proved to be sufficient for performing eigenstructure assignment. Even the knowledge of few dominant eigenpairs can be satisfactory, according to numerous papers in the literature. Moreover, eigenvalues and eigenvector can be directly obtained with appropriate measurements.

The pioneering work by Ram and Braun [10] deals exactly with the incompleteness of the original system model imposing that the assigned eigenvectors are selected within the space generated by the ones of the original system. As a consequence, the number of modifiable parameters is limited by the number of known modes. In the paper, the family of all the mathematical solutions of the modification problem is characterized. Feasibility of the solution is addressed but no algorithm to compute a feasible approximation of the theoretical solution is provided. The work of Sivan and Ram [11] is devoted to such an issue, which is overcome introducing suitable conditions that are solved in a least-square sense.

Sivan and Ram in [12] propose a method for computing mass and stiffness matrices with prescribed eigenstructure. Since the obtained mass matrix is diagonal, the method can be applied only to multiple-connected mass-spring systems. The stiffness matrix instead is computed solving an optimization problem by means of two algorithms: one that is slow but provides a global solution and another one that is more efficient but ensures the attainment of a local
minimum only. System feasibility is ensured imposing non-negative constraints. The method requires to prescribe all the modes, namely only the complete eigenstructure can be assigned.

Bucher and Braun in [13] address the eigenstructure assignment problem using modal test results only. The modified eigenvectors are forced to belong to the space spanned by the original system eigenvectors, in order to overcome the incompleteness of the eigenstructure available from the tests. Moreover, it is addressed the extraction of modal data from experimental FRFs, which is an ill-posed problem, and the sensitivity to measurement noise is examined. The method has been later improved, in fact in [14] it is proved that if both left and right eigenvectors of the system are known then it is possible to solve the eigenstructure assignment problem without truncation error.

Sivan and Ram in [15] analyse the assignment of natural frequencies with incomplete modal data. The inadequateness of the original model representation is tackled casting the assignment as an optimization problem. Feasibility of the modifications is ensured by appropriate nonnegativity constraints. It is considered also the case of interrelated modifications of mass and stiffness matrices, which is often necessary in the design of distributed-parameter systems.

## Approaches based on FRFs model

The characterization of the system through frequency response functions has several advantages, which motivate its popularity in the recent literature. Frequency response functions can be obtained from measurements, but in contrast to modal data, they are marginally affected by ill-conditioning. Therefore, such a description of the original system model is often preferable for eigenstructure assignment purposes.

For example, in [16] Mottershead managed to assign antiresonances solving the inverse eigenvalue problem for a subsidiary system. It is proved that only few measurements of the unmodified structure are required. In Kyprianou et al. [17] the assignment of natural frequencies by one added mass connected to the original system by one or two springs is studied. The problem is recast as a system of non-linear multivariate polynomial equations which are solved by means of Groebner bases. In a similar way, the same authors in [18] address the assignment of resonances or antiresonances to the original system by adding a beam, whose cross-section dimensions are to be determined.

In [19], Ouyang proposes a method for assignment of latent roots of damped asymmetric systems that consists in the modification of the mass, damping and stiffness matrices but relies on the receptance of the symmetric system. The challenge is due to the asymmetry, which causes instability. Park and Park in [20] address the eigenstructure assignment by modifying the mass, damping and stiffness matrices. The problem is recast as a linear system of equations, which can be also underdetermined (i.e. there are more modifiable parameters than equations). In [21] Mottershead and Lallement investigate the pole-zero cancellation by unit-rank modifications. In [22] Mottershead et al. deal with the same topic, focusing on the special case of repeated poles and zeros.

The general eigenstructure assignment problem is addressed by Ouyang et al. in [23], which is an extension of [8] that enable to exploit the convex optimization formulation simply from the knowledge of the FRFs. Similarly to [8], also for such a method the number of modifiable parameters is arbitrary, as well as the topology of the modification matrices and the number of assigned modes. Moreover, feasibility can be ensured specifying suitable constraints. The computed solution is globally optimal as a consequence of convexity of the formulation.

## Open issues in Dynamic Structural Modification

Although the state-of-the-art enables to solve very diverse assignment tasks by means of DSM, there are several noteworthy open problems. For example, the available methods allow assigning only complete eigenvectors, and hence require the dynamic response to be prescribed for the entire system at once (i.e. for all the degrees of freedom). In practice, however, it is often more desirable to specify exactly only the dynamic behaviour of the parts of the system that have the greatest importance. The assignment of partial eigenvectors, however, cannot be obtained with any of the methods available.

A problem that cannot be tackled satisfactorily is the possibility of modifying the dynamical properties of a system through the addition of auxiliary subsystems, despite the fact that it is of considerable practical interest, since often the original design of a system cannot be changed significantly. There are indeed several results that require the modifications of the system to have very specific topologies, but the general problem is rarely investigated.

Another relevant issue of passive control is that it is affected by spill-over. Such a phenomenon occurs when a subset of the natural frequencies is altered, thus modifying also the other frequencies. Spill-over should be avoided in order to prevent an unexpected and unwanted behaviour of the system. The available methods that enable to inhibit spill-over can deal only with simple mass-spring systems, while a method that can deal with general system models, such as the ones obtained with Finite Element modelling, is missing.

Structural modification is related also to another important problem, which is model updating, namely the tuning of a model parameters in such a way that the computed eigenvalues and eigenvectors of the model match the measured ones. Given the complexity of the Finite Element models currently employed in mechanical design, there is an increasing need of methods to perform accurate model updating, which cannot be fulfilled by the state-of-the-art.

### 1.2.2. Pole assignment through state-feedback control

Active vibration control is usually oriented to the assignment of sole eigenvalues, i.e. the socalled "pole placement". It is a very well established topic, in fact early results in such a field date back to the 1960s. For example, in [24] Wonham proved that the poles can be assigned by means of state feedback if and only if the system is controllable. In [25], Kautsky et al. investigate sensitivity of closed-loop poles with respect to perturbation of system and gain matrices. Since the pole placement problem is underdetermined for multivariable control, a technique to determine a robust solution is proposed. Ackermann in [26] introduced a wellknown formula that provide a closed form expression of control gains for single-input systems. In such a case, instead, the solution is unique, but the computation through Ackermann's formula is not numerically stable, as it requires the inversion of the controllability matrix.

The methods cited so far require to recast the second-order model (1.1) into a first-order form. In such a way, however, the matrices involved double the dimension and lose some desirable properties (e.g symmetry and sparsity). In fact, in vibration control a great effort has been put in tackling the problem directly through the second-order system model. In [27] Chu and Datta propose two methods for eigenvalue assignment and proved that if the desired poles are selfconjugated then the feedback gains are real. Datta et al. in [28] address the partial pole placement, in which chosen eigenvalues are relocated and all the other eigenvalues are left
unchanged. Orthogonality conditions for the eigenvectors of the system are derived and exploited to prove a closed form expression for the control gains.

Ram and Elhay in [29] consider the partial pole placement for multi input systems. It is proposed to solve a sequence of pole assignment tasks by single-input control and it is demonstrated that such a procedure reduces the control effort. Chu in [30] presents algorithms for robust eigenvalue assignment through state and output feedback and partial pole assignment through state feedback. In [31] Qian and Xu achieve robust partial pole placement minimizing the condition number of the closed-loop eigenvectors matrix.

The assignment of eigenstructure by active control has been first recognized by Moore in [32] and Andry et al. in [33]. Such works, that employ a first order model, characterize the class of eigenvectors that can be obtained by state feedback control. Additional references can be found in the review paper by White [34]. The assignment of partial eigenstructure using a secondorder model has been tackled by Datta et al. in [35].

Receptance-based methods have also gained popularity, since the paper by Ram and Mottershead [36] in which the assignment of poles and zeros by single-input feedback control is addressed. In [37] Tehrani et al. propose a receptance method for partial pole placement by multi-input control. The chief advantage of such methods, similarly to the case of receptancebased passive control method, is that there is no need to know the mass, damping and stiffness matrices, which is often impractical and not adequate to model the dynamic behaviour of the systems of interest.

## Open issues in active eigenstructure assignment

As pointed out in the literature review, pole placement is often oriented to the assignment of few dominant poles, which are the ones that affect the dynamics the most. The literature on partial pole placement is almost always restricted to methods for avoiding spill-over. In several cases, it is desired to perform partial pole placement allowing spill-over, as long as it can be constrained and unexpected dynamics can be prevented. In such a way, the additional freedom in pole location could be exploited to improve performances.

The assignment of eigenvectors by state-feedback control is also a relevant problem that cannot be tackled satisfactorily. In fact it can be proved that an eigenvector can be assigned if and only if it belongs to a particular vector subspace, that depends on the physical properties of the system and on the characteristics of the actuation. In particular, the dimension of such a space equals the number of actuation forces, hence for underactuated systems such a limitation can prevent the attainment of the desired eigenstructure. A design strategy that effectively deals with this issue is missing in the literature.

### 1.3. Objectives of the thesis

Although the state-of-the-art methods reviewed in the previous Section enable to achieve eigenstructure assignment in several contexts, the design of vibrating systems often requires to face issues that the state-of-the-art cannot handle. In the present thesis are proposed methods for eigenstructure assignment that aim to solve open problems emerged from the literature.

Both passive and active control have been investigated. Structural modification has been employed for the assignment of eigenvalues and partial eigenvectors in undamped systems. An approach similar to [8] has been developed, which consists in minimizing the residual of the eigenproblem equation (1.6) in a least-square sense. Given the non-convexity of the resulting optimization problem, a strategy inspired by homotopy optimization has been implemented to boost the attainment of a global minimum.

A similar optimization scheme has been applied also to the eigenstructure assignment problem through the addition of auxiliary subsystems. The main advantage of the method proposed over the state-of-the-art ones is the possibility to handle very general systems and any modification matrix which is linear in the system parameters.

The problem of eigenfrequency assignment without spill-over is also addressed. It is indeed proposed a method that enable to modify a subset of the natural frequencies of a vibrating system keeping the other ones unchanged or at least affected by very limited spill-over. Such a problem has been solved in the literature only for mass-spring systems, while the proposed method aims to avoid spill-over also for system with distributed parameters.

To conclude the contribution of the thesis to the field of Dynamic Structural Modification, a method for model updating specifically conceived for flexible link mechanisms is proposed and experimentally applied to the tuning of the Finite Element model of a planar robot.

Two issues related to active control have been also addressed. In particular, it is proposed a method for state-feedback gain synthesis that enables to perform partial pole placement through rank-1 control. The spill-over that affects the poles that are not assigned is limited employing suitable constraints. Such an approach is rarely employed in the literature and the proposed method intends to constitute a more numerically stable alternative to the available methods.

Finally, the assignment of eigenvectors by state feedback control is addressed. It is proposed a method of structural modification that enable to suitably modify the range of eigenvectors that can be obtained through state feedback control. Hence, the proposed method constitutes a hybrid active/passive strategy for eigenstructure assignment. An extensive numerical and experimental validation is also provided.

### 1.4. Outline of the thesis

The thesis is organized as follows. In Chapter 2 the assignment of eigenvalues and partial eigenvectors through passive control is investigated. In fact, eigenvectors define the distribution of vibrations over the whole system, but in several applications it would be desirable to focus on the response of the most concerning parts of the system. The proposed method allows to specify only the most relevant parts of the assigned eigenvectors, while the others are represented by optimization variables. It has been developed an optimization algorithm that is explained in detail. The proposed method is validated through two test-cases: a five-dof lumped parameter system and a 39-dof lumped-and-distributed parameter system.

The same approach can be used to tackle another relevant problem: the eigenstructure assignment through the addition of auxiliary subsystems. In Chapter 3 such a problem is addressed, employing the strategy for assignment of partial eigenvectors. In particular, the response of the primary system is assigned, while the eigenvector components correspondent to
the auxiliary systems are considered as unknowns. The proposed method is also validated with favourable results in comparison to other methods available in the literature

Chapter 4 is devoted to frequency spill-over caused by structural modification. It has been developed a novel method, that allows to inhibit spill-over, which consists of two steps: first a system with the desired eigenvalues is sought, regardless of its feasibility, and then it is computed an equivalent system which is also feasible, exploiting a special class of ordinary differential equations called isospectral flows. The method has been numerically validated with three examples: one system with lumped parameters, one with distributed parameters and finally one with lumped-and-distributed parameters.

In Chapter 5 it is presented a method for model updating specifically developed for flexible link mechanism. The proposed approach ultimately casts model updating as an optimization problem in the presence of bounds on the feasible values, by also using techniques to improve the numerical conditioning. The method is applied experimentally to the model updating of a lightweight planar robot.

Chapter 6 is devoted to a method for partial pole placement by state feedback single-input control, that allows to specify exact locations for the most important poles and a region of the complex plane to which the remaining poles are constrained to belong. To achieve such a result, a receptance based method for partial pole placement and a state-space method for regional pole placement are combined, the latter exploiting Linear Matrix Inequalities to express the regional constraint. The freedom in the placement of the non-assigned poles is exploited to minimize the controller effort. The method is validated in the control of a four-bar flexible link mechanism.

Chapter 7 investigates a limitation that affects the assignment of eigenvectors through active control. It is proposed a hybrid method that exploits passive control to modify the range of eigenvectors attainable by active control. The problem is proved to be equivalent to the minimization of the rank of a matrix. Two algorithms are proposed to compute the numerical solution. A numerical validation of the method is carried out, in order to assess the effectiveness of the proposed hybrid approach. Moreover, the two algorithms are thoroughly compared.

Chapter 8 is devoted to the experimental application of the hybrid method proposed in the previous Chapter. The employed test case is a cantilever beam made of aluminium. A modal
characterization of such a beam is performed and the Finite Element model is tuned accordingly. In order to modify the second lowest frequency mode, a piezoelectric actuator is placed on the beam. Verified that it is not possible to control adequately the beam, two suitable lumped masses are added to the beam as computed by means of the method described in the previous Chapter. Finally, the employed control scheme is described and the measurements that prove the attainment of the design specifications are provided.

## Chapter 2. Assignment of partial eigenvectors

### 2.1. Introduction

A relevant open issue in Dynamic Structural Modification is that none of the available methods can deal with design specifications that involve only some arbitrary parts of the system. Indeed they allow assigning only complete eigenvectors, and hence the dynamic response should be prescribed for the entire system at once (i.e. for all the model degrees of freedom). In contrast, it would be often desirable to specify exactly only the dynamic behaviour of the parts of the system that have the greatest importance. Hence, only some entries of the eigenvectors are exactly imposed. This problem is referred to as the partial eigenvector assignment problem and it is of great interest in particular in the presence of large scale models. Indeed, it is reasonable to expect that weakening the requests on the least concerning parts of the system might lead to a better assignment of the parts of interest, by improving the overall design of the system. Additionally, defining precise specifications for all the parts of the system, and therefore for all the degrees of freedom (dofs) of the system model employed, is often difficult. This fact may occur for instance whenever hybrid coordinates are employed (e.g. through the Craig-Bampton method)

The methods for DSM reviewed in Chapter 1, however, cannot be trivially extended to the assignment of partial eigenvectors since such a problem has some mathematical difficulties. To solve this issue, in this Chapter it is presented a novel method that allows partial eigenvector assignment, in the presence of an arbitrary number of assigned entries. The method relies on the optimization-based approach introduced in [8], and thus it inherits its peculiarities such as the possibility to assign an arbitrary number of eigenpairs, the generality of the topology of the admissible additive modifications and the ease of inclusion of feasibility constraints. The extension of method [8] to partial assignment is not straightforward, since the presence of some entries of the mode shape that are not assigned leads to a non-convex formulation of the cost function. Thus an algorithm to mitigate the issues of non-convex optimization and to boost the convergence to the global optimal solution should be adopted.

### 2.2. Problem forumlation

Let us suppose that the $N$-dof vibrating system to be modified is defined through its mass matrix $\mathbf{M} \in \mathbb{R}^{N \times N}$ and stiffness matrix $\mathbf{K} \in \mathbb{R}^{N \times N}$. Damping is neglected, as it is often done in most of the papers on structural modifications, since, in practice, several vibrating systems of interest for this kind of applications are lightly damped. These matrices represent both the original system model, in the case of improving the design and optimizing an existing system, as well as an initial guess design in the case of designing new systems.

The aim of the system modification is to find feasible additive modifications $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$ of the inertial and elastic parameters such that the modified system with mass matrix $\mathbf{M}+\Delta \mathbf{M}$ and stiffness matrix $\mathbf{K}+\Delta \mathbf{K}$ features a desired set of $N_{e} \leq N$ natural frequencies $\left(\omega_{1}, \omega_{2}, \ldots\right.$, $\omega_{N_{e}}$ ), and hence the eigenvalues, and some entries of the associated mode shapes (the eigenvectors) ( $\left.\boldsymbol{u}_{1}, \boldsymbol{u}_{2}, \ldots, \boldsymbol{u}_{N_{e}}\right)$. The topologies of $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$ are assumed a-priori in order to reflect feasible modifications of the system parameters, on the basis of technical constraints and design choices. Nevertheless, there is no limitation on the topology of these matrices, except that they should be symmetric in order to ensure the modifications to be physically meaningful. All the modifications of the system parameters, which are in practice regarded as the design variables, are then collected in vector $\boldsymbol{x} \in \mathbb{R}^{N_{x}}$. Thus it holds that $\Delta \mathbf{M}=\boldsymbol{\Delta}(\boldsymbol{x})$ and $\Delta \mathbf{K}=\Delta \mathbf{K}(\boldsymbol{x})$, and it is assumed that the dependence on $\boldsymbol{x}$ is linear. This hypothesis allows to represent a wide class of possible system modifications, as corroborated by most of the papers in literature that make the same assumption.

In order to enable the assignment of partial eigenvectors, in this work it is proposed to regard the desired eigenvectors as dependent on a set of $N_{y}<N \cdot N_{e}$ additional unknowns, denoted $\boldsymbol{y} \in \mathbb{R}^{N_{y}}$. Vector $\boldsymbol{y}$ collects all the eigenvector entries whose values are not exactly specified in the assignment, i.e. those of least interest in the design. Hence the eigenpairs to be assigned are denoted $\left(\omega_{h}, \boldsymbol{u}_{h}(\boldsymbol{y})\right)$, for $h=1, \ldots, N_{e}$.

Since both the physical modifications and the mode shape entries of the least concerning parts should be feasible, a set of constraints bounding $\boldsymbol{x}$ and $\boldsymbol{y}$ must be introduced. Such a set should include lower and upper bounds on each variable (i.e. $\boldsymbol{x}^{L} \leq \boldsymbol{x} \leq \boldsymbol{x}^{U}$, for some lower and upper bounds $\boldsymbol{x}^{L}, \boldsymbol{x}^{U} \in \mathbb{R}^{N_{x}}$, and $\boldsymbol{y}^{L} \leq \boldsymbol{y} \leq \boldsymbol{y}^{U}$, for some $\boldsymbol{y}^{L}, \boldsymbol{y}^{U} \in \mathbb{R}^{N_{y}}$, where the inequalities are element-wise). Constraints simultaneously relating different, and also inhomogeneous, variables can be also included in the design problem by means of $N_{\gamma}$ inequalities setting a convex polyhedron (defined through matrix $\mathbf{A} \in \mathbb{R}^{N_{\gamma} \chi\left(N_{x}+N_{y}\right)}$ and vector $\boldsymbol{c} \in \mathbb{R}^{N_{y}}$ ). The intersection of these two types of constraints leads to the feasible domain $\Gamma$ :

$$
\Gamma=\left\{\left\{\begin{array}{l}
\boldsymbol{x}  \tag{2.1}\\
\boldsymbol{y}
\end{array}\right\} \in \mathbb{R}^{N_{x}+N_{y}}: \boldsymbol{x}^{L} \leq \boldsymbol{x} \leq \boldsymbol{x}^{U}, \boldsymbol{y}^{L} \leq \boldsymbol{y} \leq \boldsymbol{y}^{U}, \mathbf{A}\left\{\begin{array}{l}
\boldsymbol{x} \\
\boldsymbol{y}
\end{array}\right\} \leq \boldsymbol{c}\right\}
$$

Hence, the assignment problem can be solved by finding $\left\{\begin{array}{l}x \\ y\end{array}\right\} \in \Gamma$ such that the eigenvalue problem related to each desired eigenpair is satisfied:

$$
\begin{equation*}
\left[\omega_{h}^{2}(\mathbf{M}+\Delta \mathbf{M}(x))-(\mathbf{K}+\Delta \mathbf{K}(\boldsymbol{x}))\right] \boldsymbol{u}_{h}(\boldsymbol{y})=0, \quad \text { for } h=1, \ldots, N_{e} . \tag{2.2}
\end{equation*}
$$

Since the solvability of the latter system of equations is not always ensured, in particular because of the specifications on the topology of the modification matrices and the presence of the feasible domain, the problem is recast as a constrained least-square minimization accounting simultaneously for all the desired eigenpairs:

$$
\min _{x, y}\left\{f(\boldsymbol{x}, \boldsymbol{y})=\sum_{h=1}^{N_{e}} \alpha_{h}\left\|\left[\omega_{h}^{2}(\mathbf{M}+\Delta \mathbf{M}(\boldsymbol{x}))-(\mathbf{K}+\Delta \mathbf{K}(\boldsymbol{x}))\right] \boldsymbol{u}_{h}(\boldsymbol{y})\right\|^{2},\left\{\begin{array}{l}
\boldsymbol{x}  \tag{2.3}\\
\boldsymbol{y}
\end{array}\right\} \in \Gamma\right\},
$$

The scalar coefficients $\alpha_{h}$ are introduced to reflect different levels of concern on the attainment of each eigenpair, i.e. to weigh the importance of each eigenpair. The problem can be also regularized through the approach described in [8] to boost the achievement of small norm solutions, i.e. small modifications.

It is important to remark that the cost function $f$ in problem (2.3) is, in general, non-convex. In fact, some bilinear terms of the form $x_{r} \cdot y_{s}$ (for arbitrary $r$ and $s, 1 \leq r \leq N_{x}$ and
$1 \leq s \leq N_{y}$ ), might appear in its expression, leading to a non-linear and non-convex least square problem. Therefore, a strategy to deal with this issue is needed, since standard optimization tools are likely to find local minima that are not global minima, and therefore are not reliable for achieving optimal system design. The lack of convexity of the optimization problem in Eq. (2.3) is the main difference with respect to the formulation adopted in [8], where the cost function is the squared residual of a linear system and hence is convex.

### 2.3. The proposed method

### 2.3.1. Homotopy optimization

In order to tackle the issues related to non-convex optimization, this work proposes an approach inspired by the technique of homotopy optimization (see e.g. [38, 39]). Although the attainment of a global minimal solution is not deterministically guaranteed, homotopy optimization has several advantages and has been often successfully adopted. First of all, homotopy optimization boosts effectively the convergence to the global optimal solution and therefore it improves significantly the chances to attain modifications of the system parameters leading to a very close fulfilment of the design requirements. Secondly, if compared to standard local solvers, it is less sensitive to the initial guess assumed for the numerical solution [40], once that the homotopy transformation of the cost function is performed judiciously. Additionally, compared to other global optimization methods such as for instance genetic algorithms or global search techniques, the computational effort required by homotopy is usually lower and therefore it is suitable also for medium and large scale models. In contrast, global optimization standard methods are often very time-consuming and unsuitable for large-scale and badly-conditioned systems, like eigenstructure assignment problems are. These two issues often do not allow obtaining the global optimum through such standard strategies. A novel strategy devoted to the partial eigenvector assignment problem is therefore suggested, by taking advantage of homotopy optimization.

Generally speaking, the fundamental idea behind homotopy is to replace the objective function to be minimized, $f$, with a modified one defined as an affine combination of the original
function itself and an auxiliary function, and then to solve a pattern of optimization problems. The auxiliary function, denoted $f_{c}$, is a convex simplification of the problem to be solved and its global optimum should be known or, at least, easy to be found regardless of the initial guess. The optimization starts by solving such an auxiliary convex function. Then, a morphing parameter is used to transform back the problem into the original one, while solving a sequence of optimization problems. Therefore, the last optimization problem solved is the original function $f$.

In practice, first of all the auxiliary problem $f_{c}$ is designed to approximate the original nonconvex one, $f$. Then, it is defined the new function $f_{h}$, named the homotopy map, i.e. a transformation between the auxiliary problem $f_{c}$ and the original non-convex one $f$, which depends on the scalar parameter $\lambda \in[0,1]$ :

$$
\begin{equation*}
f_{h}(\lambda)=(1-\lambda) f_{c}+\lambda f \tag{2.4}
\end{equation*}
$$

The morphing parameter $\lambda$ varies from 0 through 1 to set the transformation from $f_{c}$ to $f$. Thus, given a discretization of the interval [ 0,1 ], homotopy optimization consists in tracking the solution along such a transformation by solving a sequence of optimization problems for each value of $\lambda$. The initial guess adopted in the minimization of the non-convex function at the generic $i^{\text {th }}$ step of this sequence is the optimal solution found at the previous one. The justification of this approach is that if $\lambda$ is affected by a small variation between two consecutive iterations, it is reasonable to expect that the location of minima does not change significantly; hence a minimum can be found by means of a local solver with little effort. It has been observed in literature [41] that the final outcome of such a procedure, i.e. the minimum of $f_{h}$ with $\lambda=1$, is usually a global optimal solution of the original minimization problem $f$ if an effective convex approximation is designed. It is therefore important to carefully design both the convex auxiliary problem $f_{c}$, and the homotopy map $f_{h}$. In the next sub-sections, these two important aspects will be discussed, since an approach differing from those usually employed in literature has been proposed here.

### 2.3.2. Convex relaxation: introduction of the auxiliary variables

The convex approximation proposed in this work is achieved by replacing the bilinear terms appearing in $f$ with a suitable number $N_{t}$ of new auxiliary variables, denoted $t_{1}, \ldots, t_{N_{t}}$ and collected in vector $\boldsymbol{t} \in \mathbb{R}^{N_{t}}$. The following variable substitution is therefore performed for appropriate values of the indices $r \in\left\{1, \ldots, N_{x}\right\}, s \in\left\{1, \ldots, N_{y}\right\}$ and $i \in\left\{1, \ldots, N_{t}\right\}$ :

$$
\begin{equation*}
x_{r} \cdot y_{s} \rightarrow t_{i} \tag{2.5}
\end{equation*}
$$

The function obtained in such a way is denoted $f_{c}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t})$, and it is immediate proving its convexity. This approach is sometimes adopted in optimization and it is often referred to as "lifting", because the introduction of new variables lifts the problem into a space of higher dimension, compared with the original problem.

### 2.3.3. McCormick's constraints

Since new variables have been added, the feasible domain should be modified to include constraints on the new auxiliary variables and to represent their relation with the variables collected in $\boldsymbol{x}$ and $\boldsymbol{y}$. Ideally, since the variable $t_{i}$ replaces the bilinear term $x_{r} \cdot y_{s}$, it should be added the bilinear constraint $t_{i}-x_{r} \cdot y_{s}=0$, which is however non-convex. In order to represent this constraint by preserving convexity, an approximation is performed by taking advantage of the results demonstrated in [42], which defines the tightest convex envelopes of such constraints. This transformation is known as the so-called McCormick's relaxation [43], and can be computed given the lower and upper bounds of the variables involved. In fact, provided that $x_{r}^{L} \leq x_{r} \leq x_{r}^{U}$ and $y_{s}^{L} \leq y_{s} \leq y_{s}^{U}$, the equality $t_{i}=x_{r} \cdot y_{s}$ can be approximated by the following four inequalities:

$$
\left\{\begin{array}{l}
t_{i} \geq x_{r}^{L} \cdot y_{s}+y_{s}^{L} \cdot x_{r}-x_{r}^{L} \cdot y_{s}^{L}  \tag{2.6}\\
t_{i} \geq x_{r}^{U} \cdot y_{s}+y_{s}^{U} \cdot x_{r}-x_{r}^{U} \cdot y_{s}^{U} \\
t_{i} \leq x_{r}^{U} \cdot y_{s}+y_{s}^{L} \cdot x_{r}-x_{r}^{U} \cdot y_{s}^{L} \\
t_{i} \leq x_{r}^{L} \cdot y_{s}+y_{s}^{U} \cdot x_{r}-x_{r}^{L} \cdot y_{s}^{U}
\end{array}\right.
$$

The convex set defined through the inequalities in Eq. (2.6) for all the auxiliary variables $t_{i}$ is henceforth denoted $\Gamma_{t}$, and should be included in the feasible domain:

$$
\Gamma^{\prime}=\left\{\left\{\begin{array}{l}
\boldsymbol{x}  \tag{2.7}\\
\boldsymbol{y} \\
\boldsymbol{t}
\end{array}\right\} \in \mathbb{R}^{N_{x}+N_{y}+N_{t}}:\left\{\begin{array}{l}
\boldsymbol{x} \\
\boldsymbol{y}
\end{array}\right\} \in \Gamma, \boldsymbol{t} \in \Gamma_{t}\right\}
$$

### 2.3.4. Homotopy map definition

As for the definition of the homotopy map $f_{h}$, that morphs back the convex approximation $f_{c}$ to the original cost function $f$, a slightly different strategy is proposed in this work, compared with the typical approach which directly uses the affine combination of these two functions as the one of Eq. (2.4). In particular, by taking advantage of the problem lifting previously discussed, the homotopy map $f_{h}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, \boldsymbol{\lambda})$ is obtained from $f$ through the following substitution:

$$
\begin{equation*}
x_{r} y_{s} \rightarrow(1-\lambda) t_{i}+\lambda x_{r} y_{s} \tag{2.8}
\end{equation*}
$$

Each bilinear term $x_{r} y_{s}$ is replaced with an affine combination of itself and the associated auxiliary variable $t_{i}$, through the homotopy parameter $\lambda$. This variable transformation defines a suitable path between $f_{c}$ and $f$, since $f_{h}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, \lambda=0)=f_{c}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t})$ and $f_{h}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, \lambda=1)=f(\boldsymbol{x}, \boldsymbol{y})$, for any $\left\{\boldsymbol{x}^{\mathrm{T}} \boldsymbol{y}^{\mathrm{T}} \boldsymbol{t}^{\mathrm{T}}\right\}^{\mathrm{T}} \in \Gamma^{\prime}$ and $\left\{\boldsymbol{x}^{\mathrm{T}} \boldsymbol{y}^{\mathrm{T}}\right\}^{\mathrm{T}} \in \Gamma$. For any other $\lambda$ such that $0<\lambda<1, f_{h}$ defines a cost function in the variables $(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t})$ suitable to be employed in the intermediate iterations of homotopy optimization.

### 2.3.5. Implementation issues of the algorithm

The implementation of the proposed method deserves some additional comments. In practice, it is necessary to discretize the interval $[0,1]$ where $\lambda$ ranges and then to minimize iteratively the function $f_{h}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, \lambda)$ over $\Gamma^{\prime}$ for the selected values of $\lambda$. It is suggested to space linearly such a range through a uniform discretization step $\Delta \lambda, \Delta \lambda=1 /\left(N_{\lambda}-1\right)\left(N_{\lambda}\right.$ denotes the
number of iterations), i.e. $\lambda \in\{0, \Delta \lambda, 2 \Delta \lambda, \ldots, 1-\Delta \lambda, 1\} \subset[0,1]$. The result is not affected by the choice of $\Delta \lambda$, unless unreasonable "large values" are assumed.

The problem solved in the first iteration, i.e. the minimization of $f_{h}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, \boldsymbol{\lambda}=0)$ over $\Gamma^{\prime}$, is a convex optimization program, thus a global minimum can be found regardless of the initial guess $\left(\boldsymbol{x}_{0}, \boldsymbol{y}_{0}, \boldsymbol{t}_{0}\right)$. Afterwards, the cost function is not convex in the other iterations. Therefore, supposed that a minimizer $\left(\boldsymbol{x}_{\text {opt }}^{\lambda=\bar{\lambda}}, \boldsymbol{y}_{\text {opt }}^{\lambda=\bar{\lambda}}, \boldsymbol{t}_{\text {opt }}^{\lambda=\bar{\lambda}}\right) \in \Gamma^{\prime}$ of $f_{h}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, \bar{\lambda})$ has been found for a certain $\lambda=\bar{\lambda} \in\{0, \Delta \lambda, 2 \Delta \lambda, \ldots, 1-\Delta \lambda, 1\}$, then at the subsequent step the function $f_{h}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, \bar{\lambda}+\Delta \lambda)$ is minimized by using $\left(\boldsymbol{x}_{\text {opt }}^{\lambda=\bar{\lambda}}, \boldsymbol{y}_{\text {opt }}^{\lambda=\bar{\lambda}}, \boldsymbol{t}_{\text {opt }}^{\lambda=\bar{\lambda}}\right)$ as the initial guess. This process is continued until $\lambda=1$, i.e. when the original non-convex problem is morphed back. Finally, the optimal and feasible modifications of the mass and stiffness matrices are computed as $\Delta \mathbf{M}_{\text {opt }}=\Delta \mathbf{M}\left(\boldsymbol{x}_{\text {opt }}^{\lambda=1}\right)$ and $\Delta \mathbf{K}_{\text {opt }}=\Delta \mathbf{K}\left(\boldsymbol{x}_{\text {opt }}^{\lambda=1}\right)$.

The algorithm is briefly summarized in the pseudo code of Table 2.1.

Table 2.1. Pseudo code of the method implementation.
Input: system model $\mathbf{M}, \mathbf{K}$; admissible modifications $\Delta \mathbf{M}(\boldsymbol{x}), \Delta \mathbf{K}(\boldsymbol{x})$; feasibility constraints $\Gamma$; desired eigenpairs $\left(\omega_{h}, \boldsymbol{u}_{h}(y)\right)$; number of iterations of the homotopy map $N_{\lambda}$

Output: optimal feasible modifications $\Delta \mathbf{M}_{o p t}=\Delta \mathbf{M}\left(\boldsymbol{x}_{o p n}\right), \Delta \mathbf{K}_{o p t}=\Delta \mathbf{K}\left(\boldsymbol{x}_{o p n}\right)$

## Description

$1 \quad$ Formulation of the cost function $f$ according to Eq. (2.3)
2 Introduction of new auxiliary variables $t_{i}$
3 Definition of the homotopy map $f_{h}$ by the symbolic substitution in Eq. (2.5)

4 Symbolic computation of the McCormick's constraint and definition of $\Gamma^{\prime}$ (see Eq. (2.6) and (2.7))
$5 \quad$ Set $i:=0$ and $\lambda_{i}:=0$
6 Choose arbitrary $\left(\boldsymbol{x}_{0}, \boldsymbol{y}_{0}, \boldsymbol{t}_{0}\right)$ and solve the convex program over $\Gamma^{\prime}: \min \left\{f_{n}(\boldsymbol{x}, \boldsymbol{y}, \boldsymbol{t}, 0)\right\}$
$7 \quad$ while $i<N_{\lambda}$ do

8
Set $i:=i+1$
$9 \quad$ Set $\lambda_{i}:=\lambda_{i-1}+\Delta \lambda$

10 Solve the non-convex program $\min \left\{f_{h}\left(x, y, t, \lambda_{i}\right)\right\}$ over $\Gamma^{\prime}$ taking $\left(\boldsymbol{x}_{i-1}, \boldsymbol{y}_{i-1}, \boldsymbol{t}_{i-1}\right)$ as initial guess and set $\left(\boldsymbol{x}_{i}, \boldsymbol{y}_{i}, \boldsymbol{t}_{i}\right)$ as the optimal point
end while
12 Minimize $f(\boldsymbol{x}, \boldsymbol{y})$ over $\Gamma$ taking as initial guess $\left(\boldsymbol{x}_{N_{\lambda}-1}, \boldsymbol{y}_{N_{\lambda}-1}\right)$ and set $\left(\boldsymbol{x}_{\text {op }}, \boldsymbol{y}_{\text {op }}\right)$ as the optimal point

13
Set $\Delta \mathbf{M}_{o p}:=\Delta \mathbf{M}\left(\boldsymbol{x}_{o p}\right)$ and $\Delta \mathbf{K}_{o p}:=\Delta \mathbf{K}\left(\boldsymbol{x}_{o p}\right)$

### 2.4. Numerical examples and method validation

### 2.4.1. Five degrees of freedom lumped parameter system

In the present Section, a first example of partial assignment is examined, in order to prove the effectiveness of the proposed method. The system to which structural modification is applied is the 5-dof system shown in Figure 2.1. It consists of five lumped masses ( $m_{1}, m_{2}, m_{3}, m_{4}$ and $\left.m_{5}\right)$, connected to the frame by linear springs, whose stiffnesses are respectively $k_{g 1}, k_{g 2}, k_{g 3}$, $k_{g 4}$ and $k_{g 5}$. The springs that connect two adjacent masses are, instead, $k_{12}=75.14 \mathrm{kN} / \mathrm{m}$, $k_{23}=67.74 \mathrm{kN} / \mathrm{m}, k_{34}=75.47 \mathrm{kN} / \mathrm{m}$ and $k_{45}=83.40 \mathrm{kN} / \mathrm{m} \cdot a$

Such a system has been employed to validate other methods for eigenstructure assignment in several papers, such as [9, 23, 44]. Papers [9, 23] also provide the experimental characterization of the system modal properties.

The solution of the optimization problems has been performed through the interior point algorithm of Mosek for convex programming [45] and the interior point algorithm of fmincon from the Optimization Toolbox by Mathworks [46] for non-convex programming. The homotopy has been implemented exploiting the YALMIP [47] modelling language for Matlab.

In the test, it is supposed that dynamic response of the masses $m_{1}, m_{3}$ and $m_{5}$ is much more significant with respect to the other ones, whose eigenvector entries are just required to belong to prescribed intervals. In particular, the assignment of two natural frequencies and the associated three entries of the two eigenvectors is required. Besides showing the results


Figure 2.1. Picture of the five degrees of freedom test-case.
obtained, the comparison with the assignment of two complete eigenpairs obtained in [9] is proposed, where, in contrast, complete assignment is performed. On the one hand the numerical results are aimed at showing the capability of the novel method to compute the correct modifications by taking into account design specifications that involve only some degrees of freedom. On the other, the comparison with the complete assignment demonstrates that partial assignment boosts the achievement of the mode shape specifications at the dofs of interest.

The values of the masses and the springs of the original (unmodified) system are taken from [9] and are shown in the second column of Table 2.2. It is supposed that only the five masses and the five ground springs can be modified. Therefore the physical modifications sets the following

Table 2.2. Inertial and elastic parameters of the original and of the modified system in the first test-case.

|  | Original <br> values | Constraints | Full <br> assignment $[9]$ | Proposed <br> method |
| :--- | :--- | :--- | :--- | :--- |
| $m_{1}(\mathrm{~kg})$ | 1.727 | $[0 ;+1.08]$ | +1.0800 | +0.8829 |
| $m_{2}(\mathrm{~kg})$ | 5.123 | $[0 ;+1.08]$ | +0.1103 | +0.4977 |
| $m_{3}(\mathrm{~kg})$ | 8.214 | $[0 ;+1.08]$ | +0.9819 | +0.0013 |
| $m_{4}(\mathrm{~kg})$ | 2.609 | $[0 ;+1.08]$ | +0.0360 | +0.4784 |
| $m_{5}(\mathrm{~kg})$ | 1.339 | $[0 ;+1.08]$ | +1.0800 | +0.9841 |
| $k_{g 1}(\mathrm{kN} / \mathrm{m})$ | 94.26 | $[0 ;+144.9]$ | +80.3229 | +63.4018 |
| $k_{g 2}(\mathrm{kN} / \mathrm{m})$ | 94.26 | $[0 ;+144.9]$ | +62.5002 | +78.5887 |
| $k_{g 3}(\mathrm{kN} / \mathrm{m})$ | 94.26 | $[0 ;+144.9]$ | +144.9000 | +125.4621 |
| $k_{g 4}(\mathrm{kN} / \mathrm{m})$ | 94.26 | $[0 ;+144.9]$ | +0.0000 | +39.8018 |
| $k_{g 5}(\mathrm{kN} / \mathrm{m})$ | 94.26 | $[0 ;+144.9]$ | +106.5498 | +102.3471 |

unknown vector $\boldsymbol{x}=\left\{\begin{array}{lllllllll}\Delta m_{1} & \Delta m_{2} & \Delta m_{3} & \Delta m_{4} & \Delta m_{5} & \Delta k_{g 1} & \Delta k_{g 2} & \Delta k_{g 3} \Delta k_{g 4} \Delta k_{g 5}\end{array}\right\}^{\mathrm{T}}$. As far as the physical constraints $\boldsymbol{x}^{L}$ and $\boldsymbol{x}^{U}$ are concerned, in accordance with the cited paper, it is supposed that the mass modifications are constrained between 0 and 1.08 kg , while the stiffness modifications must be greater than 0 and less than $144.9 \mathrm{kN} / \mathrm{m}$.

The goal of the assignment is to compute the suitable parameter modifications to assign two eigenfrequencies $\omega_{1}$ and $\omega_{2}$, and two associated eigenvectors, named $\boldsymbol{u}_{1}$ and $\boldsymbol{u}_{2}$, whose first, third and fifth entries are requested to match those defined in [9]. In contrast, the second and the fourth entries are not assigned and then considered as unknowns, thus the vector of the unassigned eigenvector entries is $\boldsymbol{y}=\left\{\begin{array}{llll}u_{1,2} & u_{1,4} & u_{2,2} & u_{2,4}\end{array}\right\}^{\mathrm{T}}$ (the first subscript denotes the number of the assigned mode shape; the second one the number of the dof). The desired eigenpairs are shown in Table 2.3: it should be noticed that the normalization of the eigenvector is different from the one adopted in [9], and it is such that the first component of the first mode and the fifth component of the second mode are set to 1 . Nonetheless, the final result is not affected by the eigenvector normalization. As for the unassigned components, Table 2.3 supplies the lower and upper bounds $\boldsymbol{y}^{L}, \boldsymbol{y}^{U}$. Indeed, the corresponding eigenvector entries have been constrained with appropriate intervals including the values requested in the cited paper [8], in order to avoid unfeasible and unreasonable displacements of the second and the fourth mass. To make clearer the distinction between assigned and unassigned components, the latter ones are printed with italic font in Table 2.3.

Table 2.3. Eigenstructure comparison.

|  | Desired eigenstructure | Full <br> assignment [9] | Partial assignment | Desired eigenstructure | Full <br> assignment [9] | Partial <br> assignment |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $h=1$ |  |  | $h=2$ |  |  |
| $u_{h, 1}$ | 1.0000 | 1.0000 | 1.0000 | 0.0048 | 0.0049 | 0.0048 |
| (\% error) |  | (0\%) | (0\%) |  | (2.01\%) | (0.00\%) |
| $u_{h, 2}$ | [-0.5455;-0.1818] | -0.3626 | -0.3443 | [-0.0190;-0.0019] | -0.0098 | -0.0090 |
| $u_{h, 3}$ | 0.0545 | 0.0501 | 0.0545 | 0.0571 | 0.0589 | 0.0571 |
|  |  | (8.07\%) | (0.00\%) |  | (3.15\%) | (0.00\%) |
| $u_{h, 4}$ | [-0.0364;-0.0182] | -0.0235 | -0.0208 | [-0.9524;-0.3810] | -0.7148 | -0.6091 |
| $u_{h, 5}$ | -0.0455 | -0.0432 | -0.0455 | 1.0000 | 1.0000 | 1.0000 |
|  |  | (5.05\%) | (0.00\%) |  | (0\%) | (0\%) |
| Cosine | - | 0.9998 | 1.0000 | - | 0.9990 | 1.0000 |
| Angle [ ${ }^{\circ}$ ] | - | 1.01 | 0.00 | - | 2.53 | 0.00 |
| $\omega_{h}[\mathrm{~Hz}]$ | 50 | 49.68 | 50.00 | 60 | 60.13 | 60.00 |
| Error [Hz] | - | 0.32 | 0.00 | - | 0.13 | 0.00 |

The assignment task is also summarized in Figure 2.1, which sketches the studied system, where bold lines indicate the physical parameters whose modification is allowed, and dashed lines are used to denote the unassigned components of the eigenvectors.

The optimal modifications obtained through complete assignment and through partial assignment are listed, respectively, in the fourth and fifth column of Table 2.2. The eigenpairs of the modified systems are shown in Table 2.3. Two aggregate parameters are used to compare the outcomes of the methods: the frequency absolute error and the cosine between the desired and the obtained partial eigenvectors, computed with reference to the three desired entries. The
results clearly show that the proposed method actually provides exact assignment both on the eigenvector specifications and on the eigenvalues and satisfies all the design constraints (listed in the second and fifth columns of Table 2.3). Such result is expected because the residual of the cost function (2.3) is $5.834 \cdot 10^{-9}$. In contrast, if the complete eigenvector is forced to assume the values proposed in [8], the specifications are partially missed both in terms of frequency and mode shape, as a consequence of the constraints on the allowable modifications. Table 2.3 shows that a small adjustment in the desired response of the second and fourth mass, compared with those of the complete assignment, results in a substantial improvement in the attainment of the desired dynamic behaviour for the other three masses of interest.

The comparison of the two strategies of eigenstructure assignment clearly confirms the hypotheses that relaxing the requirements on the degrees of freedom of least concern results in a more accurate fulfilment of the specifications of interest (both in terms of mode shapes and eigenfrequencies).

As a further proof of the results obtained, Figure 2.2 shows the absolute value of some sample Frequency Response Functions (FRFs) (i.e. $H_{i, 5}(\omega), i=1, \ldots, 5$ ), obtained from $\left[\omega^{2} \mathbf{M}-\mathbf{K}\right]^{-1}$ of the original system (dashed line) and from $\left[\omega^{2}(\mathbf{M}+\Delta \mathbf{M})-(\mathbf{K}+\Delta \mathbf{K})\right]^{-1}$ of the system modified by the proposed method (solid line). A small damping has been assumed for clarity of representation. The comparison with the vertical lines, defining the desired eigenfrequencies, clearly shows the presence of the resonance peaks at the frequencies of interest for the modified system. In contrast, the original system misses the design specifications.


Figure 2.2. System FRFs.

### 2.4.2. Linear vibratory feeder

The second test case employed for numerical validation is a linear vibratory feeder, which is a system of practical relevance in the field of design of mechanical systems. Such a machinery is used to convey small products by means of the oscillatory motion of a surface, called the tray. The system model, shown in Figure 2.3, has 39 dofs, and hence this system is a challenging test to assess the method.


Figure 2.3. 39-dof model of the vibratory feeder.

Both the upper and the lower beams are 2 meters long and are modelled through 8 EulerBernoulli beam finite elements. The bending stiffnesses of these elements are $E J_{U}=4.5 \cdot 10^{3} \mathrm{Nm}^{2}$ and $E J_{L}=4.5 \cdot 10^{4} \mathrm{Nm}^{2}$ (subscripts $U$ and $L$ denote respectively the lower and the upper beam), while the linear mass densities are $\rho A_{U}=12.7 \mathrm{~kg} / \mathrm{m}$ and $\rho A_{L}=16 \mathrm{~kg} / \mathrm{m}$. Two lumped springs, whose stiffness is $k_{f 1}=k_{f 2}=50 \mathrm{~N} / \mathrm{m}$, connect the lower beam to the rigid frame. Six elastic suspensions, modelled as lumped springs, are placed at $0.25 \mathrm{~m}, 0.75 \mathrm{~m}$, $1.25 \mathrm{~m}, 1.75 \mathrm{~m}$ and also at the two ends of the beams in order to support the tray. The system is excited by three electromagnetic actuators placed at $0.5 \mathrm{~m}, 1 \mathrm{~m}$ and 1.5 m , which are modelled as concentrated masses cantilevered to the upper beam by three elastic elements. The nominal values of such parameters in the original system design are shown in Table 2.4. It is supposed that all the lumped parameters can be modified, except for the two ground springs. In contrast, none of the distributed parameters can be altered. Hence, $\boldsymbol{x}=\left\{\Delta m_{1}, \ldots, \Delta m_{3}, \Delta k_{1}, \ldots, \Delta k_{9}\right\}^{T}$.

Just one eigenpair is assigned since vibratory feeders are forced through single harmonic excitation whose frequency should match the one of a mode of vibration featuring a uniform vertical oscillation of the tray, in such a way that the products are uniformly conveyed. The vertical displacements are prescribed to equal 1 (although the method does not require any specific normalization of the desired eigenvector). The elastic rotation of the upper tray and the motion of the lower beam instead belong to the unknown vector $\boldsymbol{y}$. The displacements of the three counterweights are required to be the same, thus are represented by one variable only, named $u_{z_{C W}}$. In practice, the unknown vector is $\boldsymbol{y}=\left\{u_{\vartheta_{1}}, \ldots, u_{\vartheta_{9}}, u_{z_{10}}, \ldots, u_{z_{18}}, u_{\vartheta_{10}}, \ldots, u_{\vartheta_{18}}, u_{z_{C W}}\right\}^{T}$.

The variables $\boldsymbol{x}$ are constrained according to the second column of Table 2.4. The variables $\boldsymbol{y}$, instead, are required to belong to the interval $[-3,3]$, except $u_{z_{C W}}$ that must belong to $\left[-3,-10^{-8}\right]$ (negative values are chosen to force the counterweights to oscillate in anti-phase with respect to the upper beam. As for the eigenfrequency, it is required to be 40 Hz , which is a frequency that is often assumed for conveying small parts.

The implemented homotopy optimization method succeed in finding a global optimum of the cost function (2.3), in fact the residual is $4.806 \cdot 10^{-8}$. The computed optimal modifications of the physical parameters are shown in the third column of Table 2.4. The resulting system features a resonance at $\omega=40.000 \mathrm{~Hz}$, whose corresponding mode shape matches quite accurately the design specifications: in fact the cosine between the first 18 components of the obtained eigenvector (those related to the upper beam) and the optimal vector $(1,0,1,0, \ldots)^{T}$ (that represents, for each node, uniform vertical displacement and zero degrees rotation) is 0.9995 , thus very close to the ideal value 1 .

Table 2.4. Inertial and elastic parameters of the original and of the modified system in the second test-case.

|  | Original values | Constraints | Optimal values |
| :---: | :---: | :---: | :---: |
| $m_{1}(\mathrm{~kg})$ | 20 | $[0 ;+50]$ | +0.3070 |
| $m_{2}(\mathrm{~kg})$ | 20 | $[0 ;+50]$ | +0.9649 |
| $m_{3}(\mathrm{~kg})$ | 20 | $[0 ;+50]$ | +0.3070 |
| $k_{1}(\mathrm{kN} / \mathrm{m})$ | 8 | $[-7 ;+200]$ | +106.2 |
| $k_{2}(\mathrm{kN} / \mathrm{m})$ | 8 | $[-7 ;+200]$ | +187.5 |
| $k_{3}(\mathrm{kN} / \mathrm{m})$ | 40 | $[-30 ;+200]$ | +67.64 |
| $k_{4}(\mathrm{kN} / \mathrm{m})$ | 40 | $[-30 ;+200]$ | +67.64 |
| $k_{5}(\mathrm{kN} / \mathrm{m})$ | 8 | $[-7 ;+200]$ | +187.5 |
| $k_{6}(\mathrm{kN} / \mathrm{m})$ | 8 | $[-7 ;+200]$ | +106.2 |
| $k_{7}(\mathrm{kN} / \mathrm{m})$ | 270 | $[-200 ;+200]$ | -102.1 |
| $k_{8}(\mathrm{kN} / \mathrm{m})$ | 270 | $[-200 ;+200]$ | -96.62 |
| $k_{9}(\mathrm{kN} / \mathrm{m})$ | 270 | $[-200 ;+200]$ | -102.1 |

The obtained mode shape is depicted in Figure 2.4, where the displacement of the three concentrated masses with respect to the equilibrium is magnified five times for clarity of representation.


Figure 2.4. Mode shape at 40 Hz of the modified vibratory feeder.

In comparison, Figure 2.5 shows the mode shape of the original unmodified system with the highest participation factor at 40 Hz , whose eigenfrequency is 40.135 Hz and whose cosine, computed once again only for the first 18 components, is just 0.2395 .


Figure 2.5. Mode shape at 40.135 of the unmodified vibratory feeder.

### 2.5. Conclusion

In this Chapter it is proposed a novel method for partial eigenvector assignment through inverse dynamic structural modification. The method is suitable when the design requirements do not specify the desired mode shapes for the whole system, but rather exact requirements are made only for some parts of the system. In contrast, just upper and lower bounds are prescribed for the remaining dofs, in order to avoid unfeasible or unreasonable dynamic responses.

Such a distinction is particularly useful in two situations. First of all, when it is wanted to boost the achievement of the design specifications at some of the dofs, at the expenses of the others of least concern. Secondly, when it is not precisely known how some components of the assigned eigenvector should be specified. The latter condition may occur for parts of the system having minor importance, or whenever hybrid physical and non-physical coordinates are employed (e.g. through the Craig-Bampton method), and therefore it is difficult prescribing the dynamic response of the non-physical coordinates.

The method therefore deserves both theoretical and practical interest in the field of mechanical design. On the one hand it proposes an effective method to solve an open issue in the literature on inverse structural modification, by tackling effectively some mathematical issues. On the other, it provides a useful tool in designing and optimizing vibrating systems.

Besides solving an open issue, another novel contribution of the method is its formulation, which is based on the problem lifting to make it convex (i.e. the introduction of some auxiliary variables, constrained through a tight convex feasible domain) and then the use of the homotopy optimization to overcome the non-convex nature of the partial assignment problem and to boost the achievement of the global optimal solution.

The method has been tested numerically with two benchmark systems, often employed in the literature on inverse dynamic structural modification: a five-dof system made by lumped masses and springs and a linear vibratory feeder modeled through 39 dofs. Concerning the first testcase, the partial assignment of two eigenpairs (by specifying only three entries of each eigenvector) has been performed with very effective results. First of all, it is shown that the proposed strategy based on the homotopy optimization effectively boosts the achievement of the global optimal solution. Then, it has been assessed the initial assumption claiming that weakening the requirements on some dofs results in a more accurate assignment of the others. Indeed it is a duty of the proposed method to determine the best assignment for the dofs of least concern, in order to meet more closely the design specifications on the dofs of interest. The second test-case, which is more challenging, also allows to assess the method in a meaningful design example.

## Chapter 3. Auxiliary subsystems

### 3.1. Introduction

The most popular way to employ structural modification for the design and the optimization of vibrating systems is certainly to focus on the modification of the inertial and elastic parameters (damping is usually neglected) of the existing components of the system, without modifying the number of degrees of freedom. The possibility of modifying the dynamical properties of a system through the addition of auxiliary subsystems is rarely investigated, instead. However, such a problem is of considerable practical interest because, when the original design of a system cannot be changed, the only possible alternative is to add one or more external subsystems to the primary one.

The direct problem related to structural modifications that increase the number of degrees of freedom has been tackled in the recent literature, especially concerning distributed modifications, such as in the works by Canbaloğlu and Özgüven [48] and Hang, Shankar and Lai $[49,50]$. The inverse problem, instead, is mainly devoted to vibration suppression, as in [51] in which multiple tuned mass dampers are exploited to control a bridge-like structure, or in [52] that deals with the control of a vibrating plate by a proper choice and positioning of dynamic absorbers. The assignment of natural frequencies by added masses and springs also appeared in literature over the last years. In particular Kyprianou, Mottershead and Ouyang in [17] recast the problem as a system of non-linear multivariate polynomial equations and propose a method to solve it by means of the Groebner bases. A different insight on this approach is provided by the recent work of Liu et al. [44], in which the non-linear equations are effectively solved in a least-square sense. Moreover, they also take into account the system mode shapes, making the method suitable for eigenstructure assignment through the addition of multiple mass-spring systems.

In this Chapter, it is presented a method to perform structural modifications by means of added auxiliary systems, that allows to assign the desired dynamic behaviour of the original system both in terms of natural frequencies and mode shapes. The formulation adopted, which casts the assignment as an optimization problem, allows to attain a dynamic behaviour which is more
fitting to the desired one, in comparison to that of the original system, even when an exact assignment is not possible because of feasibility constraints. The topology of the auxiliary systems must be pre-determined but it can be chosen as arbitrary as possible. Moreover, it is possible to ensure the feasibility of the solution defining a constraint set in the optimization variable space.

### 3.2. Problem theoretical development

### 3.2.1. Problem formulation

Let us consider an undamped vibrating system, that should be modified to feature a desired set of eigenpairs. The system is modeled through its mass matrix $\mathbf{M}_{\mathbf{0}} \in \mathbb{R}^{N_{o} \times N_{o}}$ and stiffness matrix $\mathbf{K}_{\mathbf{0}} \in \mathbb{R}^{N_{o} \times N_{o}}$, where $N_{o}$ denotes the number of degrees of freedom (dofs). Damping is neglected, as it is usually assumed in the literature on structural modification, since most of the vibrating systems of interest for this kind of applications are lightly damped. The system free response is governed by the first-order eigenvalue problem:

$$
\begin{equation*}
\left[\omega_{\mathrm{o}}^{2} \mathbf{M}_{\mathbf{o}}-\mathbf{K}_{\mathbf{o}}\right] \mathbf{u}_{\mathbf{o}}=\mathbf{0} \tag{3.1}
\end{equation*}
$$

where $\omega_{0}$ is the eigenfrequency and $\mathbf{u}_{0} \in \mathbb{R}^{N_{o}}$ is the related real eigenvector (the subscript o denotes the original system).

The system modification, aimed at assigning the prescribed dynamic behaviour through the assignment of an arbitrary number of eigenpairs, is performed by adding one or more additional "external" systems having arbitrary topology and number of dofs, denoted $N_{e}$. It is only assumed that such a topology is pre-defined, in accordance with technical constraints and specifications. It is important to restate that the number of assigned eigenpairs is arbitrary and it is often less than the total number of degrees of freedom.

For ease of explanation it is also initially assumed that just one eigenpair is assigned and that the original masses and stiffnesses of the original "primary" system $\left(\mathbf{M}_{\mathbf{0}}, \mathbf{K}_{\mathbf{0}}\right)$ are not modified. Both these hypotheses will be lately relaxed and the method will be extended to account for the
simultaneous assignment of more eigenpairs, and also in the presence of modifications of the primary system.

Structural modification through external systems causes the system free response to be governed by the following $\left(N_{o}+N_{e}\right)$-dimensional eigenvalue problem:

$$
\left[\omega^{2}\left[\begin{array}{cc}
\mathbf{M}_{\mathbf{o}} & \mathbf{M}_{\mathrm{eo}}  \tag{3.2}\\
\mathbf{M}_{\mathbf{e o}}{ }^{T} & \mathbf{M}_{\mathbf{e}}
\end{array}\right]-\left[\begin{array}{cc}
\mathbf{K}_{\mathbf{o}}+\Delta \mathbf{K}_{\mathrm{o} b} & \mathbf{K}_{\mathrm{eo}} \\
\mathbf{K}_{\mathrm{eo}}{ }^{T} & \mathbf{K}_{\mathbf{e}}
\end{array}\right]\right]\left\{\begin{array}{l}
\mathbf{u}_{\mathbf{o}} \\
\mathbf{u}_{\mathbf{e}}
\end{array}\right\}=\mathbf{0}
$$

where $\mathbf{M}_{\mathbf{e}}, \mathbf{K}_{\mathbf{e}}$ are the $N_{e}$-dimensional mass and stiffness matrices of the external system, while $\mathbf{M}_{\mathrm{eo}}, \mathbf{K}_{\mathrm{eo}} \in \mathbb{R}^{N_{o} \times N_{e}}$ represent the inertial and elastic coupling matrices. The $N_{e}-$ dimensional vector $\mathbf{u}_{\mathbf{e}}$ is the portion of the eigenvector describing the motion of the external subsystem.

The aim of the system modification is to find feasible values of the entries of matrices $\mathbf{M}_{\mathbf{e}}$, $\mathbf{K}_{\mathrm{e}}, \mathbf{M}_{\mathrm{e} \boldsymbol{o}}$ and $\mathbf{K}_{\mathrm{eo}}$, as well as of the modifications of $\mathbf{K}_{\mathbf{o}}$ for the boundary dofs at the interface with the auxiliary system (denoted $\Delta \mathbf{K}_{\mathrm{ob}}$ ) such that the overall modified system in Eq. (3.2) has the desired eigenfrequency $\tilde{\omega}$, and the related eigenvector features the prescribed mode shape $\tilde{\mathbf{u}}_{\mathbf{0}} \in \mathbb{R}^{N_{o}}$ that is set for just the dofs of the primary system. No specifications are instead prescribed for the motion of the external system, whose mode shape entries $\mathbf{u}_{e}$ can assume arbitrary values, albeit feasible and reasonable.

The choice of $\tilde{\boldsymbol{\omega}}$ and $\tilde{\mathbf{u}}_{\mathbf{o}}$ depends on the specific features of the system to be optimized and of the tasks it should perform. For example, in the design of resonators it is wanted that an eigenfrequency matches the excitation frequency, and the corresponding eigenvector defines the spatial distribution of the vibrations (see e.g. [53]). Another meaningful example is vibration confinement, in which vibration at a specific frequency (e.g. those of a disturbance) is required to be confined in just a part of the system (see e.g. [54])

The topologies of the matrices of the external system are assumed a-priori to reflect the feasible layout of the external subsystem, defined on the basis of technical constraints and of the design choices. It is worth noticing that there is no limitation on the topology of these matrices, except that they must be physically meaningful. Therefore a wide range of problems can be tackled
through the proposed method. It is actually recommendable to choose such matrices as general as possible, in order to improve the chances to obtain the desired eigenstructure. The method itself will neglect the modifications that are not necessary, if there are any. Moreover, the least desirable modifications can be penalized by means of regularization techniques that will be detailed in the following.

All the modifications of the system parameters, which are in practice the design variables, are then collected in the unknown vector $\mathbf{x} \in \mathbb{R}^{N_{x}}$. Thus it holds that $\mathbf{M}_{\mathbf{e}}=\mathbf{M}_{\mathbf{e}}(\mathbf{x}), \mathbf{K}_{\mathrm{e}}=\mathbf{K}_{\mathbf{e}}(\mathbf{x})$, $\mathbf{M}_{\mathrm{eo}}=\mathbf{M}_{\mathrm{eo}}(\mathbf{x}), \mathbf{K}_{\mathrm{eo}}=\mathbf{K}_{\mathrm{eo}}(\mathbf{x})$ and $\Delta \mathbf{K}_{\mathrm{o} b}=\Delta \mathbf{K}_{\mathrm{o} b}(\mathbf{x})$, and it is assumed that the dependence on $\mathbf{x}$ is linear. This hypothesis allows representing a wide class of possible system modifications, as corroborated by most of the papers in literature assuming the same representation, since a proper definition of $\mathbf{x}$ allows also representing correlated modifications or modifications of continuous (distributed) mass and stiffness parameters (see e.g. [8]).

Even though assignment is of interest only for the eigenvector entries concerning the original subsystem, the whole system model should be accounted for in the eigenvalue problem, and therefore also $\mathbf{u}_{\mathbf{e}}$ should be included in the problem formulation. In particular, it is proposed to regard these eigenvector entries as additional unknowns in the assignment problem. The resulting unknown vector is therefore the $\left(N_{x}+N_{e}\right)$-dimensional column vector $\left\{\begin{array}{ll}\mathbf{x}^{T} & \mathbf{u}_{\mathrm{e}}{ }^{T}\end{array}\right\}^{T}$. In principle, such unknowns are not independent, in fact they are implicitly related by (3.2). Since it is difficult to make such relation explicit, unless local approximations are used, it is advantageous to regard the unknowns as independent. The numerical validation of the method has proved that such an assumption has no significant drawbacks.

Since also physical feasibility of the values of the external system parameters should be assured, a constraint set $\Gamma$ bounding $\mathbf{x}$ and $\mathbf{u}_{\mathrm{e}}$ is introduced. Such a set represents lower and upper bounds on each variable (i.e. $\mathbf{x}^{L} \leq \mathbf{x} \leq \mathbf{x}^{U}$ and $\mathbf{u}_{\mathrm{e}}{ }^{L} \leq \mathbf{u}_{\mathrm{e}} \leq \mathbf{u}_{\mathrm{e}}{ }^{U}$, given the prescribed lower and upper bounds $\mathbf{x}^{L}, \mathbf{x}^{U} \in \mathbb{R}^{N_{x}}$ and $\mathbf{u}_{\mathrm{e}}{ }^{L}, \mathbf{u}_{\mathrm{e}}{ }^{U} \in \mathbb{R}^{N_{e}}$, where the inequalities are element-wise). While the variables $\mathbf{x}$ are bounded due to technical and economical reasons, the variables $\mathbf{u}_{\mathbf{e}}$ represent the relative displacements of the auxiliary system with respect to the primary one, thus the
bounds must be set in accordance with layout constraints, safety requirements (i.e. presence of obstacles) or other design choices.

Constraints simultaneously relating different, and also inhomogeneous, variables can be also included in the design problem by means of inequalities setting a convex polyhedron (defined through proper real coefficient matrix $\mathbf{A}$ and vector $\mathbf{c}$ ). The intersection of these two convex sets of constraints leads to the feasible domain $\Gamma$, which is still a convex set:

$$
\Gamma=\left\{\left\{\begin{array}{c}
\mathbf{x}  \tag{3.3}\\
\mathbf{u}_{\mathbf{e}}
\end{array}\right\} \in \mathbb{R}^{N_{x}+N_{e}}: \mathbf{x}^{L} \leq \mathbf{x} \leq \mathbf{x}^{U}, \mathbf{u}_{\mathrm{e}}{ }^{L} \leq \mathbf{u}_{\mathrm{e}} \leq \mathbf{u}_{\mathrm{e}}{ }^{U}, \mathbf{A}\left\{\begin{array}{c}
\mathbf{x} \\
\mathbf{u}_{\mathbf{e}}
\end{array}\right\} \leq \mathbf{c}\right\}
$$

The assignment problem is therefore finding $\left\{\begin{array}{c}\mathbf{x} \\ \mathbf{u}_{\mathbf{e}}\end{array}\right\} \in \Gamma$ such that the eigenvalue problem related to the desired eigenpair is satisfied:

$$
\left\{\begin{array}{c}
{\left[\tilde{\omega}^{2} \mathbf{M}_{\mathrm{eo}}(\mathbf{x})-\mathbf{K}_{\mathrm{eo}}(\mathbf{x})\right] \mathbf{u}_{\mathrm{e}}-\Delta \mathbf{K}_{\mathrm{ob}}(\mathbf{x}) \tilde{\mathbf{u}}_{\mathrm{o}}}  \tag{3.4}\\
{\left[\tilde{\omega}^{2} \mathbf{M}_{\mathrm{eo}}(\mathbf{x})-\mathbf{K}_{\mathrm{eo}}(\mathbf{x})\right]^{T} \tilde{\mathbf{u}}_{\mathbf{o}}+\left[\tilde{\omega}^{2} \mathbf{M}_{\mathrm{e}}(\mathbf{x})-\mathbf{K}_{\mathrm{e}}(\mathbf{x})\right] \mathbf{u}_{\mathrm{e}}}
\end{array}\right\}=\left\{\begin{array}{c}
-\left[\tilde{\omega}^{2} \mathbf{M}_{\mathrm{o}}-\mathbf{K}_{\mathrm{o}}\right] \tilde{\mathbf{u}}_{\mathrm{o}} \\
\mathbf{0}
\end{array}\right\}
$$

The right-hand-side term collects the known terms, related to the desired eigenpair and to the matrices of the primary system. The left-hand-side terms collect instead all the unknowns: the terms $\left[\tilde{\omega}^{2} \mathbf{M}_{\text {eo }}(\mathbf{x})-\mathbf{K}_{\text {eo }}(\mathbf{x})\right]^{T} \tilde{\mathbf{u}}_{\mathrm{o}}$ and $\Delta \mathbf{K}_{\mathrm{ob}}(\mathbf{x}) \tilde{\mathbf{u}}_{\mathrm{o}}$ are linear functions of the unknowns and can be recast as a linear system $\mathbf{U x}=\mathbf{b}$, through a suitable matrix $\mathbf{U}$ and vector $\mathbf{b}$ (see Ref. [8]). In contrast, $\left[\tilde{\omega}^{2} \mathbf{M}_{\mathrm{eo}}(\mathbf{x})-\mathbf{K}_{\mathrm{eo}}(\mathbf{x})\right] \mathbf{u}_{\mathrm{e}}$ and $\left[\tilde{\omega}^{2} \mathbf{M}_{\mathbf{e}}(\mathbf{x})-\mathbf{K}_{\mathrm{e}}(\mathbf{x})\right] \mathbf{u}_{\mathbf{e}}$ have some non-linear terms. In particular, bilinear terms of the type $x_{i} \cdot u_{e, j}$ appear, where $x_{i}$ and $u_{e, j}$ are entries of the unknown vectors $\mathbf{x}$ and $\mathbf{u}_{\mathbf{e}}$ (for some indexes $i$ and $j$ ).

Since, in the general case, the solvability of the system of equations in Eq. (3.4) is not always assured, e.g. because of the specifications on the topology of the modification matrices and the presence of the feasible domain, the problem is recast as a constrained minimization of the norm of the residual of Eq. (3.4), to be solved through numerical optimization algorithms:

$$
\min _{\mathbf{x}, \mathbf{u}_{\mathbf{e}}}\left\{\left\|\left\{\begin{array}{c}
{\left[\tilde{\omega}^{2} \mathbf{M}_{\mathbf{e o}}(\mathbf{x})-\mathbf{K}_{\mathrm{eo}}(\mathbf{x})\right] \mathbf{u}_{\mathbf{e}}-\Delta \mathbf{K}_{\mathbf{o} b}(\mathbf{x}) \tilde{\mathbf{u}}_{\mathbf{o}}+\left[\tilde{\omega}^{2} \mathbf{M}_{\mathbf{o}}-\mathbf{K}_{\mathbf{o}}\right] \tilde{\mathbf{u}}_{\mathbf{o}}}  \tag{3.5}\\
{\left[\tilde{\omega}^{2} \mathbf{M}_{\mathbf{e o}}(\mathbf{x})-\mathbf{K}_{\mathrm{eo}}(\mathbf{x})\right]^{T} \tilde{\mathbf{u}}_{\mathbf{o}}+\left[\tilde{\omega}^{2} \mathbf{M}_{\mathbf{e}}(\mathbf{x})-\mathbf{K}_{\mathbf{e}}(\mathbf{x})\right] \mathbf{u}_{\mathbf{e}}}
\end{array}\right\}\right\|_{2}^{2},\left\{\begin{array}{c}
\mathbf{x} \\
\mathbf{u}_{\mathbf{e}}
\end{array}\right\} \in \Gamma\right\}
$$

In the case more pairs $\left(\tilde{\omega}_{i}, \tilde{\mathbf{u}}_{\mathbf{o}, i}\right)$ are required, $i=1, \ldots, N_{p}$, the problem is cast as the weighed sum of the residuals of the inverse eigenvalue problem of each eigenpair:

$$
\min _{\mathbf{x}, \mathbf{u}_{\mathbf{e}}}\left\{\sum_{i=1}^{N_{p}} \alpha_{i}\left\|\left\{\begin{array}{c}
{\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{e o}}(\mathbf{x})-\mathbf{K}_{\mathbf{e o}}(\mathbf{x})\right] \mathbf{u}_{\mathbf{e}, i}-\Delta \mathbf{K}_{\mathbf{o} b}(\mathbf{x}) \tilde{\mathbf{u}}_{\mathbf{o}, i}+\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{o}}-\mathbf{K}_{\mathbf{o}}\right] \tilde{\mathbf{u}}_{\mathbf{o}, i}}  \tag{3.6}\\
{\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{e o}}(\mathbf{x})-\mathbf{K}_{\mathbf{e o}}(\mathbf{x})\right]^{T} \tilde{\mathbf{u}}_{\mathbf{o}, i}+\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{e}}(\mathbf{x})-\mathbf{K}_{\mathbf{e}}(\mathbf{x})\right] \mathbf{u}_{\mathbf{e}, i}}
\end{array}\right\}\right\|_{2}^{2},\left\{\begin{array}{c}
\mathbf{x} \\
\mathbf{u}_{\mathbf{e}}
\end{array}\right\} \in \Gamma\right\}
$$

where $\mathbf{u}_{\mathbf{e}}=\left\{\begin{array}{lll}\mathbf{u}_{\mathbf{e}, 1}{ }^{T} & \ldots & \mathbf{u}_{\mathrm{e}, N_{p}}{ }^{T}\end{array}\right\}^{T}$ collects all the unknowns of the mode shape entries of the external system. The real, positive, scalar coefficients $\alpha_{i}$ are introduced to reflect different levels of concern on the attainment of each eigenpair. The problem can be also regularized through the approach described in [8] to improve its numerical conditioning and to boost the achievement of small norm solutions, i.e. small modifications.

### 3.2.2. Inclusion of modifications of the primary system

In the case modifications of the primary systems are allowed, through additive modification mass matrix $\Delta \mathbf{M}_{\mathbf{0}} \in \mathbb{R}^{N_{o} \times N_{o}}$ and stiffness matrix $\Delta \mathbf{K}_{\mathbf{0}} \in \mathbb{R}^{N_{o} \times N_{o}}$ (the latter collects all the modifications of the stiffness parameters, both at the interface dofs and in the internal ones), the problem extension is straightforward, due to the kind of formulation adopted. Let us, once again, assume that the topologies of these two modification matrices are a priori specified on the basis of the parameters considered to be modifiable, as it is common in practice: design requirements and constraints usually impose that the system modifications can be performed on just a reduced subset of masses and stiffnesses. By including in the unknown vector $\mathbf{x}$ also the unknown entries of $\Delta \mathbf{M}_{\mathbf{o}}$ and $\Delta \mathbf{K}_{\mathbf{0}}$, the assignment problem is written as the following function that couples all the modification matrices:

$$
\min _{\mathbf{x}, \mathbf{u}_{\mathbf{e}}}\left\{\sum_{i=1}^{N_{p}} \alpha_{i}\left\|\left\{\begin{array}{c}
{\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{e o}}(\mathbf{x})-\mathbf{K}_{\mathbf{e o}}(\mathbf{x})\right] \mathbf{u}_{\mathbf{e}, i}+\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{o}}-\mathbf{K}_{\mathbf{o}}\right] \tilde{\mathbf{u}}_{\mathbf{o}, i}+\left[\tilde{\omega}_{i}^{2} \Delta \mathbf{M}_{\mathbf{o}}(\mathbf{x})-\Delta \mathbf{K}_{\mathbf{o}}(\mathbf{x})\right] \tilde{\mathbf{u}}_{\mathbf{o}, i}}  \tag{3.7}\\
{\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{e o}}(\mathbf{x})-\mathbf{K}_{\mathbf{e o}}(\mathbf{x})\right]^{T} \tilde{\mathbf{u}}_{\mathbf{o}, i}+\left[\tilde{\omega}_{i}^{2} \mathbf{M}_{\mathbf{e}}(\mathbf{x})-\mathbf{K}_{\mathbf{e}}(\mathbf{x})\right] \mathbf{u}_{\mathbf{e}, i}}
\end{array}\right\}\right\|_{2}^{2}\right\}
$$

The newly introduced term, $\left[\tilde{\omega}_{i}^{2} \Delta \mathbf{M}_{\mathbf{o}}(\mathbf{x})-\Delta \mathbf{K}_{\mathbf{o}}(\mathbf{x})\right] \tilde{\mathbf{u}}_{\mathrm{o}, i}$, can be written as a linear function of the unknowns, and therefore Eq. (3.7) can be solved as Eq. (3.6).

### 3.3. Problem solution

The abovementioned bilinear terms lead to a non-linear least squares minimization problem in a convex domain, which is however a non-convex optimization problem. Therefore, standard optimization tools are likely to find local minima that are not global minima, and therefore are not reliable for the purpose.

In order to tackle these issues, it is proposed to exploit the homotopy-based strategy explained in the previous Chapter.

### 3.3.1. The proposed solution

The method proposed in the previous Chapter applies to the optimization of (3.7). In fact, the convex approximation $f_{c}$ is achieved replacing the bilinear terms in $f$ with a suitable number $N_{t}$ of new auxiliary variables $t_{1}, \ldots, t_{N_{t}}$ which are collected in vector $\mathbf{t} \in \mathbb{R}^{N_{t}}$ :

$$
\begin{equation*}
x_{r} u_{e, s} \rightarrow t_{i} \tag{3.8}
\end{equation*}
$$

(for appropriate values of the indices $r \in\left\{1, \ldots, N_{x}\right\}, s \in\left\{1, \ldots, N_{e}\right\}$ and $i \in\left\{1, \ldots, N_{t}\right\}$ ).

The auxiliary function obtained in this way is denoted $f_{c}\left(\mathbf{x}, \mathbf{u}_{\mathrm{e}}, \mathbf{t}\right)$.

The non-convex constraints $x_{r} u_{e, s}-t_{i}=0$ are approximated by taking advantage of the results demonstrated in [42, 43]. Provided that $x_{r}^{L} \leq x_{r} \leq x_{r}^{U}$ and $u_{e, s}^{L} \leq u_{e, s} \leq u_{e, s}^{U}$, the McCormick's inequalities hold:

$$
\left\{\begin{array}{l}
t_{i} \geq x_{r}^{L} \cdot u_{e, s}+u_{e, s}{ }^{L} \cdot x_{r}-x_{r}^{L} \cdot u_{e, s}{ }^{L}  \tag{3.9}\\
t_{i} \geq x_{r}^{U} \cdot u_{e, s}+u_{e, s}{ }^{U} \cdot x_{r}-x_{r}^{U} \cdot u_{e, S} \\
t_{i} \leq x_{r}^{U} \cdot u_{e, s}+u_{e, s}{ }^{L} \cdot x_{r}-x_{r}^{U} \cdot u_{e, s}{ }^{L} \\
t_{i} \leq x_{r}^{L} \cdot u_{e, s}+u_{e, s}{ }^{U} \cdot x_{r}-x_{r}^{L} \cdot u_{e, S}{ }^{U}
\end{array}\right.
$$

The convex set defined through Eq. (3.9), herein denoted $\Gamma_{t}$ for brevity, is included in the convex feasible domain of the auxiliary function and the homotopy map:

$$
\Gamma^{\prime}=\left\{\left\{\begin{array}{c}
\mathbf{x}  \tag{3.10}\\
\mathbf{u}_{\mathbf{e}} \\
\mathbf{t}
\end{array}\right\} \in \mathbb{R}^{N_{x}+N_{e}+N_{t}}:\left\{\begin{array}{c}
\mathbf{x} \\
\mathbf{u}_{\mathbf{e}}
\end{array}\right\} \in \Gamma\right\} \cap \Gamma_{\mathbf{t}}
$$

The homotopy map $f_{h}\left(\mathbf{x}, \mathbf{u}_{e, s}, \mathbf{t}, \lambda\right)$ is obtained from $f\left(\mathbf{x}, \mathbf{u}_{\mathbf{e}}\right)$ through the substitution of each bilinear term $x_{r} u_{e, s}$ with an affine combination of $x_{r} u_{e, s}$ itself and the associated auxiliary variable:

$$
\begin{equation*}
x_{r} u_{e, s} \rightarrow(1-\lambda) t_{i}+\lambda x_{r} u_{e, s} \tag{3.11}
\end{equation*}
$$

It is easy to verify that $f_{h}\left(\mathbf{x}, \mathbf{u}_{\mathbf{e}}, \mathbf{t}, \lambda=0\right)=f_{c}\left(\mathbf{x}, \mathbf{u}_{\mathbf{e}}, \mathbf{t}\right)$ and $f_{h}\left(\mathbf{x}, \mathbf{u}_{\mathbf{e}}, \mathbf{t}, \lambda=1\right)=f_{c}\left(\mathbf{x}, \mathbf{u}_{\mathbf{e}}\right)$, for any $\left\{\mathbf{x}, \mathbf{u}_{\mathbf{e}}, \mathbf{t}\right\} \in \Gamma^{\prime}$ and $\left\{\mathbf{x}, \mathbf{u}_{\mathbf{e}}\right\} \in \Gamma$. For any other $\lambda \in(0,1), f_{h}\left(\mathbf{x}, \mathbf{u}_{\mathbf{e}}, \mathbf{t}, \lambda\right)$ is suitable to be employed in the intermediate iterations of homotopy optimization.

The optimal modifications of the external and the primary systems are computed as those collected in the solution of the function with $\lambda=1$, denoted as $\mathbf{x}_{\text {opt }}^{\lambda=1}: \mathbf{M}_{\mathbf{e}}\left(\mathbf{x}_{\text {opt }}^{\lambda=1}\right), \mathbf{K}_{\mathbf{e}}\left(\mathbf{x}_{\text {opt }}^{\lambda=1}\right)$, $\mathbf{M}_{\mathbf{e o}}\left(\mathbf{x}_{\text {opt }}^{\lambda=1}\right), \mathbf{K}_{\mathrm{eo}}\left(\mathbf{x}_{\text {opt }}^{\lambda=1}\right), \Delta \mathbf{M}_{\mathbf{0}}\left(\mathbf{x}_{\text {opt }}^{\lambda=1}\right)$ and $\Delta \mathbf{K}_{\mathbf{0}}\left(\mathbf{x}_{\text {opt }}^{\lambda=1}\right)$.

### 3.4. Numerical examples

### 3.4.1. Example 1

In this Section an EA task is performed in two different ways, in order to prove the effectiveness and versatility of the proposed method. The primary system to which structural modification is applied is the 5-dof system shown in Figure 3.1. The system consists of five lumped masses $\left(m_{1}, m_{2}, m_{3}, m_{4}\right.$ and $\left.m_{5}\right)$, each one connected to a rigid frame by a linear spring, whose stiffnesses are respectively $k_{g 1}, k_{g 2}, k_{g 3}, k_{g 4}$ and $k_{g 5}$. Four springs also connect each mass


Figure 3.1. Original five degrees of freedom system
Table 3.1 Physical parameters of the original system of Example 1

|  | $m_{1}(\mathrm{~kg})$ | $m_{2}(\mathrm{~kg})$ | $m_{3}(\mathrm{~kg})$ | $m_{4}(\mathrm{~kg})$ | $m_{5}(\mathrm{~kg})$ |
| :--- | :---: | :---: | :---: | :---: | :---: |
| Nominal <br> value | 1.727 | 5.123 | 8.214 | 2.609 | 1.339 |
|  | $k_{g}(\mathrm{kN} / \mathrm{m})$ | $k_{12}(\mathrm{kN} / \mathrm{m})$ | $k_{23}(\mathrm{kN} / \mathrm{m})$ | $k_{34}(\mathrm{kN} / \mathrm{m})$ | $k_{45}(\mathrm{kN} / \mathrm{m})$ |
| Nominal <br> value | 94.26 | 75.14 | 67.74 | 75.47 | 83.40 |

with the contiguous ones, $k_{12}, k_{23}, k_{34}$ and $k_{45}$. The values of the masses and stiffnesses of the original primary system are taken from [9] and are shown in Table 3.1.

The homotopy map has been implemented exploiting the YALMIP [29] modeling language for Matlab. The interior point algorithm of Mosek for convex programming [30], and the interior point algorithm of fmincon from the Optimization Toolbox by Mathworks for non-convex programming have been adopted.

The external system has been assumed as made by four masses and eight lumped and massless springs, and its topology is depicted in Figure 3.2. This layout has been selected as a sample, to show the method capability to deal with different configurations. None of the original masses and springs can be instead modified. The vector of physical unknown $\mathbf{x}$ is therefore defined as $\mathbf{x}=\left\{\begin{array}{llllllllllll}m_{6} & m_{7} & m_{8} & m_{9} & k_{26} & k_{46} & k_{17} & k_{67} & k_{18} & k_{38} & k_{39} & k_{59}\end{array}\right\}^{T}$ (see Figure 3.2 for the meaning of the symbols). Each mass and spring of the external system is constrained by the bounds $\mathbf{x}^{L}$ and $\mathbf{x}^{U}$

Table 3.2. Structural modifications for Example 1

|  | Constraints | Optimal value <br> (without <br> regularization) | Optimal value <br> (with <br> regularization) |
| :---: | :---: | :---: | :---: |
| $m_{6}(\mathrm{~kg})$ | $[0 ;+8]$ | 0.2658 | 1.5874 |
| $m_{7}(\mathrm{~kg})$ | $[0 ;+8]$ | 0.1994 | 0.0000 |
| $m_{8}(\mathrm{~kg})$ | $[0 ;+8]$ | 1.3873 | 0.1331 |
| $m_{9}(\mathrm{~kg})$ | $[0 ;+8]$ | 0.1865 | 0.1765 |
| $k_{26}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 21.0401 | 38.8636 |
| $k_{46}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 0.0000 | 0.0000 |
| $k_{17}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 44.9824 | 0.0000 |
| $k_{67}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 4.6473 | 0.0000 |
| $k_{18}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 0.0000 | 26.2731 |
| $k_{38}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 50.0000 | 0.0000 |
| $k_{39}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 2.8485 | 0.9090 |
| $k_{59}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+50]$ | 50.0000 | 50.0000 |

stated in the first column of Table 3.2. Additionally, an upper bound on the overall mass is set: $m_{6}+m_{7}+m_{8}+m_{9} \leq 20 \mathrm{~kg}$. Other constraints, such as for example $k_{18}+k_{38}>0$, are imposed in order to guarantee that each auxiliary mass is actually connected to the primary system.

The two tests carried out differ in the regularization in order to represent two different design strategies. In the first case there is no regularization term to reflect the same concern about all the masses. In contrast in the second case it is supposed that the mass $m_{7}$ should be avoided if possible, and thus the penalty term $\gamma\left[\left(\frac{m_{7}}{m_{7}^{U}}\right)^{2}+\left(\frac{k_{17}}{k_{17}^{U}}\right)^{2}+\left(\frac{k_{67}}{k_{67}^{U}}\right)^{2}\right]$ is added to the cost function in Eq. (3.7) (with $\gamma$ the regularization parameter, chosen as $10^{9}$, and $m_{7}^{U}, k_{17}^{U}, k_{67}^{U}$ the upper bound of the feasible modifications, adopted to use dimensionless quantities in the penalty term).

The displacements of the five original masses are assigned by assuming the two eigenpairs also required in [9] (just a different eigenvector normalization is used, even though it does not affect the final outcome). As for the response of the four additional masses, the specified upper and lower bounds (see Table 3.3) are chosen to guarantee that the displacements of such masses do not exceed twice the size of the biggest displacement of the masses of the original system. Indeed, at the beginning of the design process there are no reasonable clues about the assignment of such additional masses, therefore the appropriate response is optimally computed by the proposed method, within a feasibility range. Hence it is defined vector $\mathbf{u}_{\mathbf{e}}= \begin{cases}u_{1,6} & u_{1,7}\end{cases}$ $\left.u_{1,8} u_{1,9} \quad u_{2,6} u_{2,7} u_{2,8} u_{2,9}\right\}$, with the meaning of the symbols explained in Figure 3.2.


Figure 3.2. Modified nine degrees of freedom system of Example 1.

The optimal modifications computed by means of the proposed homotopy method are shown in Table 3.2. The second column concerns the modifications obtained without regularization, while the third row is the case where the mass $m_{7}$ is penalized. It is worth noting that in the first case two springs, i.e. $k_{46}$ and $k_{18}$ whose optimal value is zero, should not be added to the system, whereas in the second case the parameters $m_{7}, k_{17}$ and $k_{67}$ are zero as wished, and even $k_{46}$ and $k_{38}$ are unnecessary. The obtained eigenstructure of the modified system in both cases is given in Table 3.3.

Table 3.3. Eigenstructure comparison for Example 1

|  | Original system | Desired <br> Eigenpair | Modified system (without regularization) | Modified system (with regularization) | Original system | Desired <br> Eigenpair | Modified system (without regularization) | Modified system (with regularization) |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  | $h=1$ |  |  |  | $h=2$ |  |
| $u_{h, 1}$ | 1.0000 | 1.0000 | 1.0000 | 1.0000 | 0.0028 | 0.0048 | 0.0046 | 0.0049 |
| $u_{h, 2}$ | -0.2410 | -0.3636 | -0.3640 | -0.3612 | -0.0042 | -0.0095 | -0.0094 | -0.0096 |
| $u_{h, 3}$ | 0.0259 | 0.0545 | 0.0529 | 0.0443 | 0.0341 | 0.0571 | 0.0573 | 0.0579 |
| $u_{h, 4}$ | -0.0079 | -0.0273 | -0.0156 | -0.0107 | -0.4954 | -0.7143 | -0.7249 | -0.7330 |
| $u_{h, 5}$ | -0.0205 | -0.0455 | -0.0471 | -0.0396 | 1.0000 | 1.0000 | 1.0000 | 1.0000 |
| $u_{h, \sigma}$ | - | [-2; 2] | 0.5381 | 0.1193 | - | [-2; 2] | 0.0113 | 0.0020 |
| $u_{h, 7}$ | - | [-2; 2] | 1.5853 | - | - | [-2; 2] | 0.0123 | - |
| $u_{h, \delta}$ | - | [-2; 2] | -0.0305 | 1.9983 | - | [-2; 2] | -0.0192 | 0.0178 |
| $u_{h, 9}$ | - | [-2; 2] | -0.0640 | -0.0578 | - | [-2; 2] | 1.9219 | 1.9486 |
| Cosine | 0.9930 | - | 0.9999 | 0.9998 | 0.9871 | - | 1.0000 | 0.9999 |
| Angle [ $\left.{ }^{\circ}\right]$ | 6.80 | - | 0.64 | 1.10 | 9.23 | - | 0.40 | 0.70 |
| $\omega[\mathrm{Hz}]$ | 52.44 | 50 | 50.00 | 49.98 | 64.36 | 60 | 60.27 | 60.17 |
| Error [Hz] | 2.44 | - | 0.00 | 0.02 | 4.36 | - | 0.27 | 0.17 |

The parameters used to evaluate the method effectiveness are the absolute error in the achievement of the desired eigenfrequency and the cosine of the angle between the desired eigenvector and the obtained one, computed with respect of the first 5 entries only. Such parameters clearly prove that the design specifications are met with excellent accuracy with both the configurations of the system. Indeed, in each case, both the mode shapes and the eigenfrequencies very closely approach the desired ones, with negligible errors: almost unitary cosines are obtained between the mode shapes, while the frequency errors closely approach zero. This result is remarkable, especially if considered that the parameters of the original (primary) system (to which it is wanted to assign the response) have not been altered.

As a final remark about the employed regularization technique, it is interesting to note that it has been attempted to perform the assignment without explicitly introducing the mass $m_{7}$, namely with only three additional degrees of freedom, and the obtained results are almost identical to those obtained adding the penalty terms to the cost function.

### 3.4.2. Example 2

The second example consists in a much more challenging EA task, since the assigned eigenpairs are very different from the ones of the original system. Moreover, such example allows to compare the proposed method with another state-of-the-art method [44]. The vibrating system under consideration is once again the one of the previous example, but different values of the physical parameter are assumed. Such values are shown in Table 3.4. It is supposed that the system can be modified as in Figure 3.3, namely adding five concentrated masses, each one


Figure 3.3. Modified system of Example 2.

Table 3.4 Physical parameters of the original system of Example 2

|  | $m_{1}(\mathrm{~kg})$ | $m_{2}(\mathrm{~kg})$ | $m_{3}(\mathrm{~kg})$ | $m_{4}(\mathrm{~kg})$ | $m_{5}(\mathrm{~kg})$ |
| :--- | :---: | :---: | :---: | :---: | :---: |
| Nominal <br> value | 1.73 | 5.12 | 8.21 | 2.61 | 1.34 |
|  | $k_{g}(\mathrm{kN} / \mathrm{m})$ | $k_{12}(\mathrm{kN} / \mathrm{m})$ | $k_{23}(\mathrm{kN} / \mathrm{m})$ | $k_{34}(\mathrm{kN} / \mathrm{m})$ | $k_{45}(\mathrm{kN} / \mathrm{m})$ |
| Nominal <br> value | 98.9 | 73.6 | 68.2 | 73.5 | 82.1 |

connected by an elastic element to one of the masses of the original system. Hence the vector of the unknowns is $\mathbf{x}=\left\{\begin{array}{lllllllllll}m_{6} & m_{7} & m_{8} & m_{9} & m_{10} & k_{16} & k_{27} & k_{38} & k_{49} & k_{510}\end{array}\right\}^{T}$. The assumed constraints are the same as in [44] and are shown in Table 3.5.

Two eigenpairs are assigned, respectively at $\omega_{1}=39 \mathrm{~Hz}$ and $\omega_{2}=55 \mathrm{~Hz}$. Only the components of the eigenvectors that corresponds to the original system are assigned, as in Table 3.6. The other components form the vector of unknowns $\mathbf{u}_{\mathbf{e}}=\left\{\begin{array}{llllllll}u_{1,6} & u_{1,7} & u_{1,8} & u_{1,9} & u_{1,10} & u_{2,6} & u_{2,7} & u_{2,8}\end{array} u_{2,9}\right.$ $\left.u_{2,10}\right\}$. Each variable $u_{h, i}$ is constrained within the interval $[-2,2]$, which means that the displacements of the auxiliary masses will be less than twice the size of the displacements of the masses of the original system.

It is important to note that these two eigenpairs, in contrast to those of the previous example, are very different to any eigenpair of the original system, and the specification on the presence of some nodes (see $u_{1,4}$ and $u_{2,4}$ in Table 3.6) makes the EA task challenging.

The computed optimal parameters are shown in Table 3.5, in which it is possible to compare the proposed solution to the one of [44]. The obtained eigenstructure is shown in Table 3.6, where the components of the eigenvectors relative to the auxiliary system are omitted for the sake of brevity. The absolute error on the eigenfrequencies and the cosine of the angle between the desired eigenvector and the obtained one, considered only with respect to the original system degrees of freedom, enable to assess the attainment of the specifications. It is evident that the goal is reached for both the desired eigenpairs with excellent accuracy. It is also possible to compare the results with those of [44]: the proposed method performed slightly worse in the assignment of the frequencies but slightly better in the assignment of the mode shapes.

Table 3.5 Structural modifications for Example 2

|  | Constraints | Reference [44] | Proposed method |
| :---: | :---: | :---: | :---: |
| $m_{6}(\mathrm{~kg})$ | $[0 ;+2]$ | 0.9044 | 0.9033 |
| $m_{7}(\mathrm{~kg})$ | $[0 ;+2]$ | 0.7336 | 1.0483 |
| $m_{8}(\mathrm{~kg})$ | $[0 ;+2]$ | 0.7187 | 2.0000 |
| $m_{9}(\mathrm{~kg})$ | $[0 ;+2]$ | 0.9962 | 1.2472 |
| $m_{10}(\mathrm{~kg})$ | $[0 ;+2]$ | 1.5344 | 0.2562 |
| $k_{16}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+200]$ | 108.62 | 107.8799 |
| $k_{27}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+200]$ | 93.94 | 200.0000 |
| $k_{38}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+200]$ | 73.46 | 49.2394 |
| $k_{49}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+200]$ | 61.78 | 74.8915 |
| $k_{510}(\mathrm{kN} / \mathrm{m})$ | $[0 ;+200]$ | 38.81 | 18.1657 |

Table 3.6. Eigenstructure comparison for Example 2

|  | Desired Eigenpair | Method [44] | Proposed method | Desired Eigenpair | Method [44] | Proposed method |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  | $h=1$ |  |  | $h=2$ |  |
| $u_{h, 1}$ | 1.00 | 1.000 | 1.000 | 0 | 0.000 | -0.000 |
| $u_{h, 2}$ | 0.55 | -0.538 | -0.558 | 0.01 | 0.011 | 0.009 |
| $u_{h, 3}$ | 0.20 | 0.104 | 0.224 | -0.10 | -0.278 | -0.089 |
| $u_{h, 4}$ | 0 | -0.004 | 0.000 | 0.80 | 0.864 | 0.804 |
| $u_{h, 5}$ | 0.05 | -0.002 | -0.022 | 1.00 | 1.000 | 1.000 |
| Cosine | - | 0.9955 | 0.9979 | - | 0.9910 | 1.0000 |
| Angle [ ${ }^{\circ}$ ] | - | 5.44 | 3.72 | - | 7.69 | 0.55 |
| $\omega[\mathrm{Hz}]$ | 39 | 39.00 | 39.03 | 55 | 55.01 | 54.96 |
| Error [Hz] | - | 0.004 | 0.034 | - | 0.005 | 0.038 |

### 3.5. Conclusion

In this Chapter it is proposed a novel approach to eigenstructure assignment through the use of auxiliary external systems with arbitrary topology and number of degrees of freedom. In particular, the method aims at determining optimal mass and stiffness modification parameters leading to the assignment of an arbitrary number of eigenpairs, defined in terms of eigenfrequency and mode shape for just the entries concerning the original primary system. No specifications are posed on the mode shape entries concerning the external systems, to which only upper and lower bounds are set, to avoid unfeasible and unreasonable displacements. The proposed approach casts the assignment as a minimization problem. Since some bilinear terms appear, the homotopy optimization approach of Chapter 2 has been exploited. The numerical validation has been carried out through a five-mass system, which is modified through the addition of four or five external masses (although in the first example it is proved that it is possible to lower the number of additional masses to three). The parameters of the auxiliary external system that have been computed are to alter the response of the original system in order to obtain two different eigenstructures, in both cases with good results. Excellent performances and ease of implementation are granted by the proposed method.

## Chapter 4. Passive control without spill-over

### 4.1. Introduction

The assignment of just natural frequencies (i.e. eigenvalues) is often performed through Dynamic Structural Modification. The problem that is typically addressed is the case in which only few eigenvalues are assigned, because in practice there is a small number of dominant eigenvalues that characterize the dynamic behaviour. The modification of a subset of the natural frequencies, however, can lead to undesired change of the remaining ones, which is a phenomenon known as frequency spill-over.

In active control, spill-over can be easily avoided, since there are many methods available that can perform partial eigenvalue assignment while keeping the non-assigned eigenvalues unchanged, both using a first-order $[55,56]$ and a second-order $[28,29,57]$ system model. For passive control, in contrast, few techniques for avoiding spill-over are available, and most of them are devoted to the preservation of just one specific frequency $[58,59,60]$. The general problem, instead, has been tackled only in [61], in which two methods are proposed. The first one concerns simply connected in-line mass-spring systems, while the second one deals with multiple-connected mass-spring systems.

The method proposed in this thesis intends to extend the one in [61] to general vibrating systems, with arbitrary modifiable parameters (e.g. distributed ones). Frequency assignment will be achieved in two steps: first of all, a system model that has the desired eigenvalues is sought, regardless of the feasibility constraints. After that, it is computed an equivalent system, which also complies with feasibility constraints, by integration of a matrix differential equation. The numerical issues related to the solution of the differential equation will be discussed. Finally, the proposed method is validated in three different test-cases of interest.

### 4.2. Method description

Let us suppose that the undamped vibrating system under consideration can be represented by means of the mass matrix $\mathbf{M}_{0} \in \mathbb{R}^{n \times n}$ and the stiffness matrix $\mathbf{K}_{0} \in \mathbb{R}^{n \times n}$. The system will be denoted as the pair $\left(\mathbf{M}_{0}, \mathbf{K}_{0}\right)$. It is supposed that a subset of the eigenvalues of the system, i.e. $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{p}$, is known. It is wanted to modify the system in such a way that the $p$ abovementioned eigenvalues are replaced by $\mu_{1}, \mu_{2}, \ldots, \mu_{p}$, while the other $n-p$ eigenvalues are left unchanged, if possible, or at least the shift of such eigenvalues is kept to a minimum. It is also of fundamental importance that such modifications are feasible, and that the resulting matrices are physically meaningful.

In order to achieve such a result, the proposed method relies on a two-steps scheme, similar to the one that has been proposed in [61]. First of all, it is sought a system $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ that features the desired eigenvalues, without taking into account the physical constraints, by taking advantage of the method proposed in [62]. The second step, instead, consists in the search of a system $\left(\mathbf{M}_{\infty}, \mathbf{K}_{\infty}\right)$, within the space of systems that have the same eigenvalues as $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$, such that it is as close as possible to the space of feasible systems. In order to perform the second step, a class of matrix differential equations, known as isospectral flows, has been exploited. Finally, the feasible system that provides the best approximation of the desired eigenvalues is computed as the orthogonal projection of $\left(\mathbf{M}_{\infty}, \mathbf{K}_{\infty}\right)$ onto the space of feasible systems. Figure 4.1 schematically represent the proposed method. The present Section is devoted to give a detailed explanation of the stages that constitutes the method.


Figure 4.1. Method description

### 4.2.1. Direct updating of the system

The method that is described in this Chapter has been proposed by Yang et al. in [62], thus it will be briefly explained, and for the details it is suggested to refer to the original paper.

In order to employ such a method, it is required the knowledge of the eigenvectors $\boldsymbol{u}_{1}, \ldots, \boldsymbol{u}_{p}$ correspondent to the eigenvalues to be modified $\lambda_{1}, \ldots, \lambda_{p}$. Such eigenvectors will be retained in the modified system design, after the first step, and they will be associated to the desired eigenvalues $\mu_{1}, \ldots, \mu_{p}$. This feature is not of interest, because the assignment of eigenvectors is not addressed, and in fact in the second step eigenvectors will be neglected.

For convenience, the following matrices are defined:

$$
\begin{align*}
& \Lambda_{1}=\operatorname{diag}\left(\lambda_{1}, \ldots, \lambda_{p}\right) \\
& \Sigma_{1}=\operatorname{diag}\left(\mu_{1}, \ldots, \mu_{p}\right)  \tag{4.1}\\
& \mathbf{U}_{1}=\left[\boldsymbol{u}_{1}, \ldots, \boldsymbol{u}_{p}\right]
\end{align*}
$$

It is wanted to find a new system $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ that features the same eigenpairs as $\left(\mathbf{M}_{0}, \mathbf{K}_{0}\right)$, except for $\left(\lambda_{1}, \boldsymbol{u}_{1}\right), \ldots,\left(\lambda_{p}, \boldsymbol{u}_{p}\right)$ which are to be replaced with $\left(\mu_{1}, \boldsymbol{u}_{1}\right), \ldots,\left(\mu_{p}, \boldsymbol{u}_{p}\right)$. Such a request can be split into two constraints:

1) a condition that imposes the new eigenpairs, which can be expressed in matrix form as

$$
\begin{equation*}
\left(\mathbf{M}_{0}+\Delta \mathbf{M}\right) \mathbf{U}_{1} \boldsymbol{\Sigma}_{1}=\left(\mathbf{K}_{0}+\Delta \mathbf{K}\right) \mathbf{U}_{1} \tag{4.2}
\end{equation*}
$$

2) a condition that assures that spill-over does not take place, while assigning the desired eigenpairs with Eq. (4.2).

It is not trivial to express such a condition without the knowledge of the eigenpairs to be retained. In [63], however, it has been proposed the following equation

$$
\begin{equation*}
\Delta \mathbf{K}\left(\mathbf{M}_{0}^{-1}-\mathbf{U}_{1} \mathbf{U}_{1}^{T}\right)-\Delta \mathbf{M}\left(\mathbf{M}_{0}^{-1} \mathbf{K}_{0} \mathbf{M}_{0}^{-1}-\mathbf{U}_{1} \Lambda_{1} \mathbf{U}_{1}^{T}\right)=\mathbf{0}, \tag{4.3}
\end{equation*}
$$

which requires only the knowledge of the eigenpairs to be modified, but assures that the remaining eigenpairs are left unchanged by the additive modification $\mathbf{\Delta M}$ and $\mathbf{\Delta K}$.

In order to avoid the inversion of the mass matrix $\mathbf{M}_{0}^{-1}$, in [62] it is also suggested to choose the modification matrices as

$$
\begin{align*}
\mathbf{\Delta M} & =\mathbf{M}_{0} \boldsymbol{\Psi} \mathbf{M}_{0} \\
\Delta \mathbf{K} & =\mathbf{M}_{0} \boldsymbol{\Phi} \mathbf{M}_{0} \tag{4.4}
\end{align*}
$$

where $\boldsymbol{\Psi}, \boldsymbol{\Phi} \in \mathbb{R}^{n \times n}$ are real symmetric matrices. In fact, in such a way, the two conditions (4.2) and (4.3) are equivalent to

$$
\begin{align*}
& \left(\mathbf{M}_{0}+\mathbf{M}_{0} \Psi \mathbf{M}_{0}\right) \mathbf{U}_{1} \Sigma_{1}=\left(\mathbf{K}_{0}+\mathbf{M}_{0} \Phi \mathbf{M}_{0}\right) \mathbf{U}_{1} \\
& \Phi\left(\mathbf{M}_{0}-\mathbf{M}_{0} \mathbf{U}_{1} \mathbf{U}_{1}^{T} \mathbf{M}_{0}\right)-\Psi\left(\mathbf{K}_{0}-\mathbf{M}_{0} \mathbf{U}_{1} \boldsymbol{\Lambda}_{1} \mathbf{U}_{1}^{T} \mathbf{M}_{0}\right)=\mathbf{0} . \tag{4.5}
\end{align*}
$$

It is recommendable to solve such equations for the optimal $\boldsymbol{\Phi}$ and $\Psi$ in a least-square sense. In particular, it is suggested to solve the following optimization problem:

$$
\begin{align*}
& \min _{\Phi, \Psi}\left\|\left(\mathbf{M}_{0}+\mathbf{M}_{0} \Psi \mathbf{M}_{0}\right) \mathbf{U}_{1} \boldsymbol{\Sigma}_{1}-\left(\mathbf{K}_{0}+\mathbf{M}_{0} \boldsymbol{\Phi} \mathbf{M}_{0}\right) \mathbf{U}_{1}\right\|^{2}  \tag{4.6}\\
& \quad \text { subject to } \boldsymbol{\Phi}\left(\mathbf{M}_{0}-\mathbf{M}_{0} \mathbf{U}_{1} \mathbf{U}_{1}^{T} \mathbf{M}_{0}\right)-\boldsymbol{\Psi}\left(\mathbf{K}_{0}-\mathbf{M}_{0} \mathbf{U}_{1} \boldsymbol{\Lambda}_{1} \mathbf{U}_{1}^{T} \mathbf{M}_{0}\right)=\mathbf{0} .
\end{align*}
$$

Given the convexity of the proposed formulation, it is possible to find a numerical solution with moderately low effort. In fact, convex optimization is a well understood topic in numerical computing, the state-of-the-art algorithm are highly satisfactory and there are many software packages available for the purpose. Supposed that $\hat{\Phi}$ and $\hat{\Psi}$ are the optimal matrices of problem (4.6), then the modified system can be obtained as

$$
\begin{align*}
& \mathbf{M}_{s}=\mathbf{M}_{0}+\mathbf{M}_{0} \hat{\Psi} \mathbf{M}_{0} \\
& \mathbf{K}_{s}=\mathbf{K}_{0}+\mathbf{M}_{0} \hat{\mathbf{\Phi}} \mathbf{M}_{0} . \tag{4.7}
\end{align*}
$$

An important issue is that the mass and stiffness matrices should be positive definite and positive semidefinite, respectively. The numerical validation of the method showed that the obtained matrices comply with the positive definite constraint, without explicitly requiring so. If it is necessary, however, it is possible to exploit one of the popular algorithm for semidefinite optimization, that allows to take into account constraints on the positive definiteness of the matrices. Software that implement such algorithms are, for example, Mosek [45], SeDuMi [64] and SDPT3 [65]. For the sake of efficiency, however, semidefinite constraint should be imposed only if the optimization of (4.6) fails.

### 4.2.2. Isospectral flow

The system $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ obtained in the first step meets the requirements regarding the eigenvalues, but does not correspond to a physically meaningful system. In fact, the matrices $\mathbf{M}_{s}$ and $\mathbf{K}_{s}$ are dense and unstructured matrices. Such an issue is taken into consideration in the second step.

In practice, it is wanted to find a system $(\mathbf{M}, \mathbf{K})$ that has the same eigenvalues as $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$, namely that belongs to the space

$$
\begin{equation*}
\mathcal{M}\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)=\left\{\left(\mathbf{P}^{T} \mathbf{M}_{s} \mathbf{P}, \mathbf{P}^{T} \mathbf{K}_{s} \mathbf{P}\right) \text { such that } \mathbf{P} \text { is invertible }\right\} \subseteq \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times n} \tag{4.8}
\end{equation*}
$$

Moreover, it is wanted that $(\mathbf{M}, \mathbf{K})$ belongs also to the space $\mathcal{V}$ of the feasible systems. The definition of such a subspace requires knowledge of the system design and it depends on the physical parameters that can be reasonably modified. Moreover, it is supposed that the mass
and stiffness matrices depend linearly on the modifiable parameters. Such an assumption is made in many other works about structural modification (e.g. [23]) since it allows to represent most of the types of modifications, so it does not invalidate the generality of the method.

Thus, the space $\mathcal{V} \subseteq \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times n}$ can be defined as follows:

$$
\begin{equation*}
\mathcal{V}=\left\{(\mathbf{M}, \mathbf{K}): \mathbf{M}=\mathbf{M}_{0}+\sum_{i=1}^{n_{M}} x_{i} \mathbf{M}_{i}, \mathbf{K}=\mathbf{K}_{0}+\sum_{i=1}^{n_{K}} y_{i} \mathbf{K}_{i} \text { for some coefficients } x_{i} \text { and } y_{i}\right\}, \tag{4.9}
\end{equation*}
$$

where $n_{M}$ and $n_{K}$ are, respectively, the number of mass elements and stiffness elements that can be modified. The matrices $\mathbf{M}_{i}$ and $\mathbf{K}_{i}$ represent the topology of the contribution of each modifiable parameter, while the coefficients $x_{i}$ and $y_{i}$ represent the magnitude of the modification.

From a geometrical point of view, the search of a system that complies with both the constraints on the natural frequencies and the feasibility constraints is equivalent to the search of the intersection of $\mathcal{M}\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ and $\mathcal{V}$. If such an intersection is empty, it is still interesting to find the system that minimizes the distance between such spaces, because it could be interpreted as a least-squares solution of the optimal design problem.

In order to describe the proposed method, let us define the projection $\mathcal{P}: \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times n} \rightarrow \mathcal{V}$ of a generic system $(\mathbf{X}, \mathbf{Y})$ as

$$
\begin{equation*}
\mathcal{P}(\mathbf{X}, \mathbf{Y})=\left(\mathcal{P}_{1}(\mathbf{X}), \mathcal{P}_{2}(\mathbf{Y})\right) \tag{4.10}
\end{equation*}
$$

where $\mathcal{P}_{1}$ and $\mathcal{P}_{2}$ are the orthogonal projection of the mass and stiffness matrices, respectively. The computation of the projection is very straightforward, if the matrices are considered as vectors. As an example, the algorithm for the projection of the mass matrix $\mathcal{P}_{1}$ is described. It is wanted to find the coefficients $x_{1}, \ldots, x_{n_{M}}$ such that

$$
\begin{equation*}
\mathcal{P}_{1}(\mathbf{X})=\mathbf{M}_{0}+x_{1} \mathbf{M}_{1}+\ldots+x_{n_{M}} \mathbf{M}_{n_{M}} \tag{4.11}
\end{equation*}
$$

Let us denote vec () the operator that stacks the columns of a matrix one below the other. For convenience it is also defined

$$
\begin{align*}
& \mathbf{S}_{M}=\left[\operatorname{vec}\left(\mathbf{M}_{1}\right) \ldots \operatorname{vec}\left(\mathbf{M}_{n_{M}}\right)\right]  \tag{4.12}\\
& \boldsymbol{\delta} \boldsymbol{m}=\operatorname{vec}\left(\mathbf{X}-\mathbf{M}_{0}\right)
\end{align*}
$$

By exploiting the orthogonality of the least-squares residual to the range of $\mathbf{S}_{M}$, the appropriate coefficients in (4.11) are computed solving the following problem:

$$
\begin{equation*}
\min _{x}\left\|\mathbf{S}_{M} \boldsymbol{x}-\boldsymbol{\delta} \boldsymbol{m}\right\|^{2} \tag{4.13}
\end{equation*}
$$

where $\boldsymbol{x}=\left(x_{1}, \ldots, x_{n_{M}}\right)^{T}$. Being $\hat{\boldsymbol{x}}$ the least-squares solution, the projection itself is

$$
\begin{equation*}
\mathcal{P}_{1}(\mathbf{X})=\mathbf{M}_{0}+\operatorname{vec}^{-1}\left(\mathbf{S}_{M} \hat{\boldsymbol{x}}\right) . \tag{4.14}
\end{equation*}
$$

The definition of $\mathcal{P}_{2}$ is analogous and thus it is omitted.

The projection onto $\mathcal{V}$ is the tool that allows to express the distance of a general system on $\mathcal{M}\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ to the space $\mathcal{V}$. It is worth to remark that, given the definition (4.8), any system in $\mathcal{M}\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ is uniquely determined by an invertible matrix $\mathbf{P}$, hence the distance function is dependent on $\mathbf{P}$. For clarity of notation, it is defined

$$
\begin{align*}
& \gamma_{1}(\mathbf{P})=\mathbf{P}^{T} \mathbf{M}_{s} \mathbf{P}-\mathcal{P}_{1}\left(\mathbf{P}^{T} \mathbf{M}_{s} \mathbf{P}\right) \\
& \gamma_{2}(\mathbf{P})=\mathbf{P}^{T} \mathbf{K}_{s} \mathbf{P}-\mathcal{P}_{2}\left(\mathbf{P}^{T} \mathbf{K}_{s} \mathbf{P}\right) \tag{4.15}
\end{align*}
$$

Hence, the distance function $F$ is expressed by

$$
\begin{equation*}
F(\mathbf{P})=\frac{1}{2}\left(\left\|\gamma_{1}(\mathbf{P})\right\|^{2}+\left\|\gamma_{2}(\mathbf{P})\right\|^{2}\right) \tag{4.16}
\end{equation*}
$$

Such a function has been also proposed in [66] to solve the inverse eigenvalue problem for linear pencils. In practice, the function is minimized by integrating the so-called descent flow, which is

$$
\left\{\begin{array}{l}
\dot{\mathbf{P}}(t)=-\nabla F(\mathbf{P}(t))  \tag{4.17}\\
\mathbf{P}(0)=\mathbf{I},
\end{array}\right.
$$

where $\mathbf{I}$ is the identity matrix. In the cited paper it is also given a closed form for the gradient of $F$, particularly useful for numerical computations, that is

$$
\begin{equation*}
\nabla F(\mathbf{P})=2\left(\mathbf{M}_{s} \mathbf{P} \gamma_{1}(\mathbf{P})+\mathbf{K}_{s} \mathbf{P} \gamma_{2}(\mathbf{P})\right) . \tag{4.18}
\end{equation*}
$$

A remarkable property of the descent flow is that its solution $\mathbf{P}(t)$ asymptotically converges. Thus, defined $\mathbf{P}_{\infty}=\lim _{t \rightarrow+\infty} \mathbf{P}(t)$, the system that (locally) minimizes the distance to $\mathcal{V}$ is

$$
\begin{align*}
& \mathbf{M}_{\infty}=\mathbf{P}_{\infty}^{T} \mathbf{M}_{s} \mathbf{P}_{\infty} \\
& \mathbf{K}_{\infty}=\mathbf{P}_{\infty}^{T} \mathbf{K}_{s} \mathbf{P}_{\infty} . \tag{4.19}
\end{align*}
$$

Then, the system in $\mathcal{V}$ that minimizes the distance to $\mathcal{M}\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ is exactly the projection of $\left(\mathbf{M}_{\infty}, \mathbf{K}_{\infty}\right)$, namely

$$
\begin{align*}
& \mathbf{M}=\mathcal{P}_{1}\left(\mathbf{M}_{\infty}\right) \\
& \mathbf{K}=\mathcal{P}_{2}\left(\mathbf{K}_{\infty}\right) . \tag{4.20}
\end{align*}
$$

The system ( $\mathbf{M}, \mathbf{K}$ ) computed by means of the described procedure actually correspond to a physically meaningful system, and it has the desired eigenvalues, or at least is affected by limited spill-over, as proven in Section 4.3 where different tests are shown. Despite the computed solution representing an actual mechanical system, it might not be sufficient for many design applications, in which other technical constraints have a fundamental role. Moreover, the algorithm for computing the solution of (4.17) requires a more detailed explanation if constraints are defined.

In the next Sections the constrained problem and its numerical integration are addressed.

### 4.2.3. Projections with constraints

It is evident that the unbounded space $\mathcal{V}$ also contains systems that can be hardly implemented in practice, due to technical reasons. Hence, it is recommendable to introduce constraints on the
coefficients $x_{i}$ and $y_{i}$ that appear in the definition (4.9). Collecting the coefficients into vectors $\boldsymbol{x}=\left(x_{1}, \ldots, x_{n_{M}}\right)^{T}$ and $\boldsymbol{y}=\left(y_{1}, \ldots, y_{n_{K}}\right)^{T}$, the constraint sets are defined by

$$
\begin{align*}
& \Gamma_{x}=\left\{\boldsymbol{x} \in \mathbb{R}^{n_{\mu}}: \mathbf{A}_{x} \boldsymbol{x} \leq \boldsymbol{b}_{x}, \mathbf{A}_{x, \text { eq }} \boldsymbol{x}=\boldsymbol{b}_{x, \text { eq }}, \boldsymbol{x}_{L} \leq \boldsymbol{x} \leq \boldsymbol{x}_{U}\right\} \\
& \Gamma_{y}=\left\{\boldsymbol{y} \in \mathbb{R}^{n_{K}}: \mathbf{A}_{y} \boldsymbol{y} \leq \boldsymbol{b}_{y}, \mathbf{A}_{y, \text { eq }} \boldsymbol{y}=\boldsymbol{b}_{y, \text { eq }}, \boldsymbol{y}_{L} \leq \boldsymbol{y} \leq \boldsymbol{y}_{U}\right\} \tag{4.21}
\end{align*}
$$

for some appropriate matrices $\mathbf{A}_{x}, \mathbf{A}_{x, \text { eq }}, \mathbf{A}_{y}, \mathbf{A}_{y, \text { eq }}$ and vectors $\boldsymbol{b}_{x}, \boldsymbol{b}_{x, \text { eq }}, \boldsymbol{x}_{L}, \boldsymbol{x}_{U}, \boldsymbol{b}_{y}, \boldsymbol{b}_{y, \text { eq }}$, $\boldsymbol{y}_{L}, \boldsymbol{y}_{U}$. Hence the constrained feasible space is defined by

$$
\begin{equation*}
\hat{\mathcal{V}}=\left\{(\mathbf{M}, \mathbf{K}): \mathbf{M}=\mathbf{M}_{0}+\sum_{i=1}^{n_{M}} x_{i} \mathbf{M}_{i}, \mathbf{K}=\mathbf{K}_{0}+\sum_{i=1}^{n_{K}} y_{i} \mathbf{K}_{i}, \boldsymbol{x} \in \Gamma_{x}, \boldsymbol{y} \in \Gamma_{y}\right\} . \tag{4.22}
\end{equation*}
$$

From the theoretical point of view, the projection onto $\hat{\mathcal{V}}$ is not different from the projection described in the previous Section, except for the fact that the unconstrained least-squares problem (4.13) must be replaced with a constrained least-squares problem. From the numerical point of view, in contrast, the presence of constraints considerably affects the computational performances. Considering that the projections $\mathcal{P}_{1}$ and $\mathcal{P}_{2}$ must be evaluated at each iteration of the numerical integration of (4.17), it is obvious that it is critical to be able to compute leastsquares solutions fast and accurately. Nevertheless, while testing the proposed method, it has been observed that the linear least-squares solvers implemented in commercial software gave unsatisfactory results.

Such a relevant issue has been tackled employing a scaling strategy, which allows a substantial improvement of the obtained results in terms of accuracy of the computed solution. Scaling is aimed at recasting the problem with dimensionless unknowns, which is advantageous in many ways, including the numerical point of view [67]. In practice, it is suggested to choose the basis of $\hat{\mathcal{V}}$ in such a way that all the coefficients $x_{1}, \ldots, x_{n_{M}}$ (or $y_{1}, \ldots, y_{n_{K}}$ ) have roughly order of magnitude $10^{\circ}$. This can be achieved simply multiplying by suitable scalar factors the matrices $\mathbf{M}_{i}$ and $\mathbf{K}_{i}$. In fact, given non-zero scalars $\eta_{1}, \ldots, \eta_{n_{M}}$ and $\tau_{1}, \ldots, \tau_{n_{K}}$, the space $\hat{\mathcal{V}}$ can be represented as

$$
\begin{equation*}
\hat{\mathcal{V}}=\left\{(\mathbf{M}, \mathbf{K}): \mathbf{M}=\mathbf{M}_{0}+\sum_{i=1}^{n_{M}} x_{i} \eta_{i} \mathbf{M}_{i}, \mathbf{K}=\mathbf{K}_{0}+\sum_{i=1}^{n_{K}} y_{i} \tau_{i} \mathbf{K}_{i}, \boldsymbol{x} \in \hat{\Gamma}_{x}, \boldsymbol{y} \in \hat{\Gamma}_{y}\right\} \tag{4.23}
\end{equation*}
$$

It is clear that this is equivalent to (4.22), provided that the constraint sets $\Gamma_{x}$ and $\Gamma_{y}$ are replaced properly by new constraint sets $\hat{\Gamma}_{x}$ and $\hat{\Gamma}_{y}$. A judicious choice of the $\eta_{i}$ and of the $\tau_{i}$ allows to keep the components of $\boldsymbol{x}$ and $\boldsymbol{y}$ roughly of the same magnitude, which is beneficial to the accuracy of the projection and to the CPU time.

### 4.2.4. Numerical integration

Although (4.17) can be compactly represented as one matrix differential equation, it is actually a system of $n^{2}$ scalar differential equations and thus it can be big even for small-scale system models. For such a reason, the numerical integration must be performed with care, in order to reach convergence in a reasonable time even with limited hardware resources.

The tests carried out for validating the proposed method have shown that the differential equation under consideration is stiff, thus the standard algorithms for numerical integration tend to require small discretization steps. Such an occurrence is common when the differential equation contains terms that are very different in magnitude and hence it is ill-conditioned, as it is the case in (4.18). In fact, mass and stiffness matrices can have very different orders of magnitude, in real mechanical vibrating systems.

To overcome such a problem, it is beneficial using specific algorithms for stiff equations, which are already implemented in many software packages for scientific computing. Moreover, it is suggested to employ a scaling strategy, similar to the one explained in Section 4.2.3. In practice, it is solved an equivalent problem in which the mass and stiffness matrices are scaled to match the desired order of magnitude, and also the feasible space $\mathcal{V}$ is modified accordingly. In fact, given the scaling factors $\alpha>0$ and $\beta>0$, the system $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ is replaced by the scaled system

$$
\begin{gather*}
\tilde{\mathbf{M}}_{s}=\alpha \mathbf{M}_{s} \\
\tilde{\mathbf{K}}_{s}=\beta \mathbf{K}_{s} . \tag{4.24}
\end{gather*}
$$

The space of feasible systems (4.9), instead, is replaced by:

$$
\begin{equation*}
\tilde{\mathcal{V}}=\left\{(\mathbf{M}, \mathbf{K}): \mathbf{M}=\alpha \mathbf{M}_{0}+\sum_{i=1}^{n_{M}} \tilde{x}_{i} \mathbf{M}_{i}, \mathbf{K}=\beta \mathbf{K}_{0}+\sum_{i=1}^{n_{K}} \tilde{y}_{i} \mathbf{K}_{i} \text { for some coefficients } \tilde{x}_{i} \text { and } \tilde{y}_{i}\right\} . \tag{4.25}
\end{equation*}
$$

It is easy to prove that a matrix $\mathbf{P}$ that minimizes the distance to the feasible space for the scaled system, also works for the original system. In fact, let us suppose that $\left(\mathbf{P}^{T} \tilde{\mathbf{M}}_{s} \mathbf{P}, \mathbf{P}^{T} \tilde{\mathbf{K}}_{s} \mathbf{P}\right) \in \tilde{\mathcal{V}}$, then there exist coefficients $\tilde{x}_{1}, \ldots, \tilde{x}_{n_{M}}$ such that

$$
\begin{gather*}
\mathbf{P}^{T} \tilde{\mathbf{M}}_{s} \mathbf{P}=\alpha \mathbf{M}_{0}+\sum_{i=1}^{n_{M}} \tilde{x}_{i} \mathbf{M}_{i} \\
\Leftrightarrow \quad \alpha \mathbf{P}^{T} \mathbf{M}_{s} \mathbf{P}=\alpha\left(\mathbf{M}_{0}+\sum_{i=1}^{n_{M}} \frac{\tilde{x}_{i}}{\alpha} \mathbf{M}_{i}\right) . \tag{4.26}
\end{gather*}
$$

Hence for $x_{i}=\tilde{x}_{i} / \alpha$ it holds $\mathbf{P}^{T} \mathbf{M}_{s} \mathbf{P}=\mathbf{M}_{0}+\sum_{i=1}^{n_{M}} x_{i} \mathbf{M}_{i}$. The same argument can be used to prove an analogous assertion for stiffness matrices, therefore $\left(\mathbf{P}^{T} \mathbf{M}_{s} \mathbf{P}, \mathbf{P}^{T} \mathbf{K}_{s} \mathbf{P}\right) \in \mathcal{V}$ and then the original and the scaled problems are equivalent.

It should be remarked that the scaled problem can be solved also in the presence of constraints, as described in Section 4.2.3. It is necessary, however, to compensate the effect of the factors $\alpha$ and $\beta$ by an appropriate choice of the coefficients $\eta_{i}$ and $\tau_{i}$, and also of the constraint sets $\Gamma_{x}$ and $\Gamma_{y}$. For example, if the coefficient $\eta_{i}$ is employed for a variable of the original problem, then $\hat{\eta}_{i}=\alpha \cdot \eta_{i}$ would be a recommendable choice for the same coefficient in the scaled problem.

According to the numerical validation that has been carried out, the proposed method enables to find a stationary point of Eq. (4.17) within a reasonable time span, for models that can be actually used to represent real mechanical systems with good accuracy. Some examples are shown in the next Section.

### 4.3. Numerical examples

The proposed method has been implement in Matlab, for the sake of numerical assessment. Three examples are shown in this Section, to demonstrate the capabilities of the developed strategy for assignment of natural frequencies.

It has been taken advantage of available software, whose performances are comparable to the state-of-the-art. In particular, for the optimization problem (4.6) it has been used Gurobi [68], together with the useful modelling language Yalmip [47]. For what concerns the numerical integration of the differential equation (4.17), instead, the Matlab built-in function ode15s for stiff equations has been used. In practice, since the limit for $t \rightarrow+\infty$ of the flow $\mathbf{P}(t)$ is sought, the differential equation is integrated in the interval $\left[0, T_{f}\right]$. For $T_{f} \gg 0$, the proposed procedure allows to compute a matrix which approximates the stationary point of (4.17) up to machine precision. Finally, the constrained linear least-squares problem that occur in the computation of the projections $\mathcal{P}_{1}$ and $\mathcal{P}_{2}$ have been solved with the function lsqlin included in the Optimization Toolbox for Matlab by Mathworks.

### 4.3.1. Cantilever beam

The first example is the cantilever beam pictured in Figure 4.2. The model consists of three Euler-Bernoulli beam elements of the same length.


Figure 4.2. Cantilever beam
The considered system already appeared in the literature [62,69] and hence it is assumed as a benchmark. The beam length is supposed to be $L=3 m$, hence the length of each finite element is $l=1 \mathrm{~m}$. The other model parameters, i.e. the linear mass density $\rho A$ and the product of the Young's modulus and the area moment of inertia $E J$, are shown in Table 4.1. The resulting mass and stiffness matrices are

$$
\begin{align*}
& \mathbf{M}_{0}=\left[\begin{array}{cccccc}
1.56 & 0.66 & 0.54 & -0.39 & 0 & 0 \\
0.66 & 0.36 & 0.39 & -0.27 & 0 & 0 \\
0.54 & 0.39 & 3.12 & 0 & 0.54 & -0.39 \\
-0.39 & -0.27 & 0 & 0.72 & 0.39 & -0.27 \\
0 & 0 & 0.54 & 0.39 & 3.12 & 0 \\
0 & 0 & -0.39 & -0.27 & 0 & 0.72
\end{array}\right] \\
& \mathbf{K}_{0}=\left[\begin{array}{cccccc}
12 & 18 & -12 & 18 & 0 & 0 \\
18 & 36 & -18 & 18 & 0 & 0 \\
-12 & -18 & 24 & 0 & -12 & 18 \\
18 & 18 & 0 & 72 & -18 & 18 \\
0 & 0 & -12 & -18 & 24 & 0 \\
0 & 0 & 18 & 18 & 0 & 72
\end{array}\right] \tag{4.27}
\end{align*}
$$

The eigenvalues of the original system are shown in the first column of Table 4.2. The three lowest eigenvalues are wanted to be equal to $\mu_{1}=0.05, \mu_{2}=1.50$ and $\mu_{3}=11.0$, while the others should remain the same. The system matrices computed after direct updating are

$$
\begin{align*}
& \mathbf{M}_{s}=\left[\begin{array}{ccccccc}
1.5442 & 0.6523 & 0.5324 & -0.3788 & 0.0352 & -0.0004 \\
0.6523 & 0.3567 & 0.3891 & -0.2662 & 0.0079 & -0.0007 \\
0.5324 & 0.3891 & 3.1192 & -0.0119 & 0.4651 & -0.3905 \\
-0.3788 & -0.2662 & -0.0119 & 0.7145 & 0.3798 & -0.2664 \\
0.0352 & 0.0079 & 0.4651 & 0.3798 & 3.1274 & 0.0183 \\
-0.0004 & -0.0007 & -0.3905 & -0.2664 & 0.0183 & 0.7200
\end{array}\right] \\
& \mathbf{K}_{s}=\left[\begin{array}{cccccc}
12.0203 & 18.0091 & -11.9892 & 17.9917 & -0.0121 & 0.0004 \\
18.0091 & 36.0039 & -17.9947 & 17.9971 & -0.0029 & -0.0001 \\
-11.9892 & -17.9947 & 24.0277 & 0.0037 & -11.9658 & 17.9955 \\
17.9917 & 17.9971 & 0.0037 & 72.0027 & -17.9915 & 17.9994 \\
-0.0121 & -0.0029 & -11.9658 & -17.9915 & 24.0202 & -0.0084 \\
0.0004 & -0.0001 & 17.9955 & 17.9994 & -0.0084 & 71.9992
\end{array}\right] \tag{4.28}
\end{align*}
$$

The eigenvalues of $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ are shown in the second column of Table 4.2. It is supposed that the system parameters that can be modified are the linear mass densities $\rho A_{i}$ and the products of the Young's modulus and area moment of inertia $E J_{i}$, for each beam element $i=1,2,3$. Hence the space of feasible systems is defined as in (4.9), with $n_{M}=n_{K}=3$ and

$$
\begin{aligned}
& \mathbf{M}_{1}=\frac{l}{420}\left[\begin{array}{cccccc}
156 & 22 l & 54 & -13 l & 0 & 0 \\
22 l & 4 l^{2} & 13 l & -3 l^{2} & 0 & 0 \\
54 & 13 l & 156 & -22 l & 0 & 0 \\
-13 l & -3 l^{2} & -22 l & 4 l^{2} & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0
\end{array}\right] \quad \mathbf{K}_{1}=\frac{1}{l^{3}}\left[\begin{array}{ccccccc}
12 & 6 l & -12 & 6 l & 0 & 0 \\
6 l & 4 l^{2} & -6 l & 2 l^{2} & 0 & 0 \\
-12 & -6 l & 12 & -6 l & 0 & 0 \\
6 l & 2 l^{2} & -6 l & 4 l^{2} & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0
\end{array}\right] \\
& \mathbf{M}_{2}=\frac{l}{420}\left[\begin{array}{cccccc}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 156 & 22 l & 54 & -13 l \\
0 & 0 & 22 l & 4 l^{2} & 13 l & -3 l^{2} \\
0 & 0 & 54 & 13 l & 156 & -22 l \\
0 & 0 & -13 l & -3 l^{2} & -22 l & 4 l^{2}
\end{array}\right] \mathbf{K}_{2}=\frac{1}{l^{3}}\left[\begin{array}{cccc}
0 & 0 & 0 & 0 \\
0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 12 & 6 l \\
0 & 0 & 6 l & 4 l^{2} \\
-12 & -6 l & 2 l^{2} \\
0 & 0 & -12 & -6 l \\
0 & 12 & -6 l \\
0 & 0 & 6 l & 2 l^{2} \\
-6 l & 4 l^{2}
\end{array}\right] \\
& \mathbf{M}_{3}=\frac{l}{420}\left[\begin{array}{llllll}
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 156 & 22 l \\
0 & 0 & 0 & 0 & 22 l & 4 l^{2}
\end{array}\right]
\end{aligned}
$$

Moreover, the constraints in the second column of Table 4.1 are imposed. Hence the differential equation (4.17) is integrated until $T_{f}=100$. Finally, the resulting system is

$$
\begin{align*}
& \mathbf{M}=\left[\begin{array}{cccccc}
1.56 & 0.66 & 0.54 & -0.39 & 0 & 0 \\
0.66 & 0.36 & 0.39 & -0.27 & 0 & 0 \\
0.54 & 0.39 & 3.12 & 0 & 0.54 & -0.39 \\
-0.39 & -0.27 & 0 & 0.72 & 0.39 & -0.27 \\
0 & 0 & 0.54 & 0.39 & 4.0162 & 0.3792 \\
0 & 0 & -0.39 & -0.27 & 0.3792 & 0.9268
\end{array}\right] ;  \tag{4.29}\\
& \mathbf{K}=\left[\begin{array}{cccccc}
12.3657 & 18.5485 & -12.3657 & 18.5485 & 0 & 0 \\
18.5485 & 37.0970 & -18.5485 & 18.5485 & 0 & 0 \\
-12.3657 & -18.5485 & 24.6472 & -0.1262 & -12.2815 & 18.4223 \\
18.5485 & 18.5485 & -0.1262 & 73.9416 & -18.4223 & 18.4223 \\
0 & 0 & -12.2815 & -18.4223 & 28.4567 & 5.8404 \\
0 & 0 & 18.4223 & 18.4223 & 5.8404 & 85.3700
\end{array}\right] .
\end{align*}
$$

Such a system can be actually obtained modifying the original system as in the third column of Table 4.2. The eigenvalues of the system $(\mathbf{M}, \mathbf{K})$ are shown in Table 4.2. Such results prove that the three lowest eigenvalues are actually closer to the objective, with respect to the original system design. The spill-over that affects the remaining eigenvalues, however, is very limited, in fact the relative error is about $2 \%$ for all the three cases.

Table 4.1. Computed modifications

| Parameter | Original <br> value | Constraint | Modification |
| :--- | :--- | :--- | :--- |
| $\rho A_{1}[\mathrm{~kg} / \mathrm{m}]$ | 1.4 | $[0,+\infty)$ | +0.0000 |
| $\rho A_{2}[\mathrm{~kg} / \mathrm{m}]$ | 1.4 | $[0,+\infty)$ | +0.0000 |
| $\rho A_{3}[\mathrm{~kg} / \mathrm{m}]$ | 1.4 | $[0,+\infty)$ | +0.8043 |
| $E J_{1}\left[\mathrm{Nm}^{2}\right]$ | 27 | $[0,+\infty)$ | +0.8227 |
| $E J_{2}\left[\mathrm{Nm}^{2}\right]$ | 27 | $[0,+\infty)$ | +0.6334 |
| $E J_{2}\left[\mathrm{Nm}^{2}\right]$ | 27 | $[0,+\infty)$ | +9.3940 |

Table 4.2. Obtained eigenvalues

|  | $\left(\mathbf{M}_{0}, \mathbf{K}_{0}\right)$ | $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ | $(\mathbf{M}, \mathrm{K})$ |
| :---: | :---: | :---: | :---: |
|  | $3.6346 \cdot 10^{-2}$ | $5.0000 \cdot 10^{-2}$ | $4.6117 \cdot 10^{2}$ |
|  | $1.4365 \cdot 10^{0}$ | $1.5000 \cdot 10^{0}$ | $1.5004 \cdot 10^{0}$ |
|  | $1.1470 \cdot 10^{1}$ | $1.1000 \cdot 10^{1}$ | $1.1021 \cdot 10^{1}$ |
|  | $5.8167 \cdot 10^{1}$ | $5.8167 \cdot 10^{1}$ | $5.6863 \cdot 10^{1}$ |
|  | $2.0602 \cdot 10^{2}$ | $2.0602 \cdot 10^{2}$ | $2.0151 \cdot 10^{2}$ |
|  | $8.1884 \cdot 10^{2}$ | $8.1884 \cdot 10^{2}$ | $8.3430 \cdot 10^{2}$ |
| 0000000 | 27.3082 | 0.0000 | 7.7666 |
|  | 4.2302 | 0.0000 | 0.0276 |
|  | 4.2702 | 0.0000 | 0.1927 |
|  | - | 0.0000 | 2.2407 |
|  | - | 0.0000 | 2.1915 |
|  | - | 0.0000 | 1.8883 |

### 4.3.2. Lumped parameter systems

Although for mass-spring systems the eigenvalues assignment without spill-over problem has been solved in [61], the proposed method has been tested with one of the examples in the cited paper. In particular, it is considered the 10 degrees of freedom system in Figure 4.3. The nominal value of the mass and stiffness parameters are shown in the first column of Table 4.3.


Figure 4.3. Lumped parameter system

The two lowest eigenvalues are wanted to be equal to $\mu_{1}=9012$ and $\mu_{2}=12118$, while the others should be left unchanged [61]. It is supposed, in accordance with the benchmark, that all the masses and springs can be modified. The differential equation (4.17) is solved until $T_{f}=30$. To improve the numerical integration, the scaling factors $\eta=1$ and $\tau=10^{-5}$ are employed. The computed modifications, which complies with the prescribed constraints, are shown in Table 4.3. The eigenvalues comparison, shown in Table 4.4, proves that the proposed method can assign the two desired frequencies without spill-over. Such result is comparable to the one obtained in [61], in terms of accuracy of the assignment.

Table 4.3. Computed modifications.

| Parameter | Original <br> value | Constraint | Modification |
| :---: | :---: | :---: | :---: |
| $m_{1}[\mathrm{~kg}]$ | 30 | $[-30,+\infty)$ | -1.3354 |
| $m_{2}[\mathrm{~kg}]$ | 35 | $[-35,+\infty)$ | -1.9799 |
| $m_{3}[\mathrm{~kg}]$ | 40 | $[-40,+\infty)$ | -2.6579 |
| $m_{4}[\mathrm{~kg}]$ | 45 | $[-45,+\infty)$ | -3.4365 |
| $m_{5}[\mathrm{~kg}]$ | 45 | $[-45,+\infty)$ | -4.2918 |
| $m_{6}[\mathrm{~kg}]$ | 45 | $[-45,+\infty)$ | -5.4105 |
| $m_{7}[\mathrm{~kg}]$ | 40 | $[-40,+\infty)$ | -2.1045 |
| $m_{8}[\mathrm{~kg}]$ | 35 | $[-35,+\infty)$ | -1.9893 |
| $m_{9}[\mathrm{~kg}]$ | 30 | $[-30,+\infty)$ | -1.5407 |
| $m_{10}[\mathrm{~kg}]$ | 25 | $[-25,+\infty)$ | -0.1986 |
| $k_{1}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $3.0118 \cdot 10^{4}$ |
| $k_{2}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $6.1159 \cdot 10^{4}$ |
| $k_{3}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $1.2533 \cdot 10^{5}$ |
| $k_{4}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $1.8437 \cdot 10^{5}$ |
| $k_{5}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-5.1242 \cdot 10^{4}$ |
| $k_{6}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $1.1121 \cdot 10^{5}$ |
| $k_{7}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $2.3010 \cdot 10^{5}$ |
| $k_{8}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $6.5613 \cdot 10^{3}$ |
| $k_{9}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $5.9681 \cdot 10^{4}$ |
| $k_{10}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $7.5688 \cdot 10^{4}$ |
| $k_{11}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-5.6379 \cdot 10^{4}$ |
| $k_{12}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-4.1956 \cdot 10^{4}$ |
| $k_{13}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-2.6967 \cdot 10^{4}$ |
| $k_{14}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-3.3854 \cdot 10^{4}$ |
| $k_{15}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-3.7155 \cdot 10^{4}$ |
| $k_{16}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-7.0636 \cdot 10^{4}$ |
| $k_{17}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-9.1187 \cdot 10^{3}$ |
| $k_{18}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-3.4443 \cdot 10^{4}$ |
| $k_{19}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-1.2269 \cdot 10^{5}$ |
| $k_{20}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-6.2758 \cdot 10^{4}$ |
| $k_{21}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-3.7457 \cdot 10^{4}$ |
| $k_{12}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-3.1824 \cdot 10^{4}$ |
| $k_{23}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-1.3914 \cdot 10^{4}$ |
| $k_{24}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-1.5011 \cdot 10^{4}$ |
| $k_{25}[\mathrm{~N} / \mathrm{m}]$ | $2.4 \cdot 10^{5}$ | $\left[-2.4 \cdot 10^{5},+\infty\right)$ | $-3.1275 \cdot 10^{4}$ |
|  |  |  |  |
|  |  | 2 |  |

Table 4.4. Eigenvalues comparison

|  | $\left(\mathbf{M}_{0}, \mathbf{K}_{0}\right)$ | $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ | (M,K) |
| :---: | :---: | :---: | :---: |
|  | $6.4024 \cdot 10^{3}$ | $9.0120 \cdot 10^{3}$ | $9.0120 \cdot 10^{3}$ |
|  | $9.6317 \cdot 10^{3}$ | $1.2118 \cdot 10^{4}$ | $1.2118 \cdot 10^{4}$ |
|  | $1.4622 \cdot 10^{4}$ | $1.4622 \cdot 10^{4}$ | $1.4622 \cdot 10^{4}$ |
|  | $2.2122 \cdot 10^{4}$ | $2.2122 \cdot 10^{4}$ | $2.2122 \cdot 10^{4}$ |
|  | $2.3065 \cdot 10^{4}$ | $2.3065 \cdot 10^{4}$ | $2.3065 \cdot 10^{4}$ |
|  | $2.8474 \cdot 10^{4}$ | $2.8474 \cdot 10^{4}$ | $2.8474 \cdot 10^{4}$ |
|  | $3.2227 \cdot 10^{4}$ | $3.2227 \cdot 10^{4}$ | $3.2227 \cdot 10^{4}$ |
|  | $3.5759 \cdot 10^{4}$ | $3.5759 \cdot 10^{4}$ | $3.5759 \cdot 10^{4}$ |
|  | $4.2302 \cdot 10^{4}$ | $4.2302 \cdot 10^{4}$ | $4.2302 \cdot 10^{4}$ |
|  | $4.9129 \cdot 10^{4}$ | $4.9129 \cdot 10^{4}$ | $4.9129 \cdot 10^{4}$ |
| $\begin{gathered} 00 \\ 0 \\ 0 \\ 0 \\ 0 \\ 0 \\ 0 \end{gathered}$ | 28.9570 | 0.0000 | 0.0000 |
|  | 20.5172 | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |
|  | - | 0.0000 | 0.0000 |

### 4.3.3. Linear vibratory feeder

The third example of frequency assignment is performed on a model of a linear vibratory feeder, such as the ones that are employed for product conveyance. The system, represented in Figure 4.4, consists of a beam of length $L=2 m$, modelled as four Euler-Bernoulli elements, with two elastic supports (i.e. $k_{1}$ and $k_{2}$ ) and three electromagnetic actuators, modelled as lumped masses (i.e. $m_{1}, m_{2}$ and $m_{3}$ ). Such a lumped-and-distributed parameter system, which is already present


Figure 4.4. Linear vibratory feeder
in the literature [8], has 13 degrees of freedom. The nominal values of the physical quantities that characterize the system are shown in the first column of Table 4.5. It is wanted that the third and the fourth natural frequencies match $30 \mathrm{~Hz}\left(\mu_{1}=35530.58\right)$ and 45 Hz ( $\mu_{2}=79943.80$ ), respectively, while the other frequencies should not be affected by the modifications. It is supposed that the linear mass density $\rho A$ and the product of Young's modulus and area moment of inertia $E J$ of the beam can be modified, as well as the masses $m_{i}$, for $i=1,2,3$ and the springs $k_{i}$, for $i=1, \ldots, 5$. The differential equation (4.17) is solved until $T_{f}=100$, employing the scaling factors $\eta=1$ and $\tau=10^{-7}$ to improve the numerical integration. The computed optimal modifications are shown in the third column of Table 4.5. The eigenvalues comparison in Table 4.6 enables to evaluate the obtained assignment: the difference between the desired frequencies and the obtained ones is significantly reduced and the spill-over that affects the remaining frequencies is very limited (less than $2 \%$ ). Since the
employed test-case represents an industrial machine, the proposed method is proved to constitute a proper tool for design of vibrating systems.

Table 4.5. Computed modifications

| Parameter | Original value | Constraint | Modification |
| :--- | :--- | :--- | :---: |
| $m_{1}[\mathrm{~kg}]$ | 25 | $[0,+\infty)$ | 0.0000 |
| $m_{2}[\mathrm{~kg}]$ | 25 | $[0,+\infty)$ | 0.0000 |
| $m_{3}[\mathrm{~kg}]$ | 25 | $[0,+\infty)$ | 0.0000 |
| $\rho A[\mathrm{~kg} / \mathrm{m}]$ | 15.8 | $[0,+\infty)$ | 0.0788 |
| $k_{1}[\mathrm{~N} / \mathrm{m}]$ | $1.0 \cdot 10^{3}$ | $[0,+\infty)$ | $0.0000 \cdot 10^{0}$ |
| $k_{2}[\mathrm{~N} / \mathrm{m}]$ | $1.0 \cdot 10^{3}$ | $[0,+\infty)$ | $0.0000 \cdot 10^{0}$ |
| $k_{3}[\mathrm{~N} / \mathrm{m}]$ | $8.7 \cdot 10^{5}$ | $[0,+\infty)$ | $1.7891 \cdot 10^{4}$ |
| $k_{4}[\mathrm{~N} / \mathrm{m}]$ | $8.7 \cdot 10^{5}$ | $[0,+\infty)$ | $2.6975 \cdot 10^{4}$ |
| $k_{5}[\mathrm{~N} / \mathrm{m}]$ | $8.7 \cdot 10^{5}$ | $[0,+\infty)$ | $1.7890 \cdot 10^{4}$ |
| $E J\left[\mathrm{Nm}{ }^{2}\right]$ | $4.5 \cdot 10^{5}$ | $[0,+\infty)$ | $2.1931 \cdot 10^{3}$ |

Table 4.6. Obtained eigenvalues

|  | $\left(\mathbf{M}_{0}, \mathbf{K}_{0}\right)$ | $\left(\mathbf{M}_{s}, \mathbf{K}_{s}\right)$ | (M, K) |
| :---: | :---: | :---: | :---: |
|  | $1.8747 \cdot 10^{1}$ | $1.8747 \cdot 10^{1}$ | $1.8719 \cdot 10^{1}$ |
|  | $8.6704 \cdot 10^{1}$ | $8.6704 \cdot 10^{1}$ | $8.6509 \cdot 10^{1}$ |
|  | $3.3855 \cdot 10^{4}$ | $3.5531 \cdot 10^{4}$ | $3.4766 \cdot 10^{4}$ |
|  | $7.5322 \cdot 10^{4}$ | $7.9944 \cdot 10^{4}$ | $7.6647 \cdot 10^{4}$ |
|  | $1.1401 \cdot 10^{5}$ | $1.1401 \cdot 10^{5}$ | $1.1623 \cdot 10^{5}$ |
|  | $9.4036 \cdot 10^{5}$ | $9.4036 \cdot 10^{5}$ | $9.4156 \cdot 10^{5}$ |
|  | $6.9338 \cdot 10^{6}$ | $6.9338 \cdot 10^{6}$ | $6.9342 \cdot 10^{6}$ |
|  | $2.6582 \cdot 10^{7}$ | $2.6582 \cdot 10^{7}$ | $2.6582 \cdot 10^{7}$ |
|  | $8.8765 \cdot 10^{7}$ | $8.8765 \cdot 10^{7}$ | $8.8755 \cdot 10^{7}$ |
|  | $2.1765 \cdot 10^{8}$ | $2.1765 \cdot 10^{8}$ | $2.1763 \cdot 10^{8}$ |
|  | $5.2345 \cdot 10^{8}$ | $5.2345 \cdot 10^{8}$ | $5.2339 \cdot 10^{8}$ |
|  | $1.4514 \cdot 10^{9}$ | $1.4514 \cdot 10^{9}$ | $1.4512 \cdot 10^{9}$ |
|  | $1.7691 \cdot 10^{9}$ | $1.7691 \cdot 10^{9}$ | $1.7689 \cdot 10^{9}$ |
| $\begin{aligned} & \frac{0}{0} \\ & 0 \\ & 0 \\ & 0 \\ & \dot{0} \\ & \dot{\sim} \end{aligned}$ | - | 0.0000 | 0.1627 |
|  | - | 0.0000 | 0.1372 |
|  | 4.7162 | 0.0000 | 2.1278 |
|  | 5.7813 | 0.0000 | 4.0072 |
|  | - | 0.0000 | 1.9615 |
|  | - | 0.0000 | 0.2551 |
|  | - | 0.0000 | 0.0336 |
|  | - | 0.0000 | 0.0237 |
|  | - | 0.0000 | 0.0023 |
|  | - | 0.0000 | 0.0204 |
|  | - | 0.0000 | 0.0516 |
|  | - | 0.0000 | 0.0913 |
|  | - | 0.0000 | 0.0544 |

### 4.4. Conclusion

The proposed method intends to fill a gap in the state-of-the-art, since there was no way to assign natural frequencies while avoiding spill-over, for general vibrating systems. The tests carried out show that the two-steps approach is a practical way to alter a subset of the natural frequencies by computing suitable modifications of the system parameters. Spill-over can be often prevented, or it is very moderate, whenever it occurs. The numerical difficulties arisen in the validation of the method have been effectively overcome, improving its reliability. In conclusion, the proposed method actually widens the capabilities of passive control, making partial eigenvalues assignment attainable without spill-over, even for complex systems.

The proposed method can be also employed for model updating, which consists in tuning the physical parameters of a system to make the natural frequencies of a model match those obtained from experimental measurements. An example of such an application can be found in Chapter 8.

## Chapter 5. Model updating for flexible link mechanisms

### 5.1. Introduction

The use of lightweight multibody systems, such as mechanisms or manipulators, has been increasing significantly over the last years [70]. Indeed this kind of construction ensures smaller material and power consumption, thus allowing for smaller actuators. On the other hand, the flexibility of the links makes the motion planning and control critical and hence imposes the use of advanced techniques that account for the system resonant dynamics. Therefore, great attention has been paid to dynamic modelling of flexible-link multibody systems (FLMSs) [71], in order to boost the development of effective model-based design or control techniques or to allow for numerical simulations.

Several modelling approaches are proposed in literature. A widespread approach is to take advantage of finite element (FE) methods to discretize the link flexibility through a finite number of elastic degrees of freedom (dofs) and to represent the total motion of the system as the superposition of a large amplitude rigid-body motion and the coupled small-amplitude elastic deformation [72, 73, 74].

The availability of a correct model formulation is however not sufficient to predict correctly the dynamic response. Indeed, these FE models are not effective if the inertial and elastic properties of the model are not correctly tuned. Therefore, it is necessary paying attention to the correct estimation of the model parameters, either through direct physical measurement or using experimental parameter identification techniques. The latter process of correcting the FE model parameters to feature a set of experimental measurements is known in the field of structural dynamics as model updating [75].

Model updating can be performed by following two different types of approach: non-iterative or iterative techniques [75, 76]. Non-iterative techniques (also denoted direct or one-step) provide analytical solutions to model updating problem through a single step. Nonetheless, such solutions are often without physical meaning, and scarcely robust to measurement noise.

Conversely, iterative techniques (also denoted parametric techniques) preserve the physical meaning of the updated parameters at the cost of an increase of the computational complexity due to the numerical iterations adopted to reduce the mismatch between experimental and model responses. Basically, iterative techniques compute the updated parameters through an objective function that represents the difference between analytical and experimental results. Therefore, they may lead to local optimal solutions if the updating problem is formulated as a non-convex optimization problem. In the last years, the use of iterative techniques has been boosted by the advanced techniques and toolbox for numerical optimization [77, 78].

Although, several model updating techniques have been developed and implemented in the field of structural mechanics, in the multibody field the issue of model updating has been just marginally addressed, while great attention has been paid to the development of theoretical models. This fact is mainly due to the challenges in the problem formulation for multibody systems. Indeed, on the one hand, dynamic models representing FLMSs are nonlinear. On the other hand, flexible-link multibody models are typically formulated by means of elastic and rigid coordinates, which simplify the mathematical development of the FE model [74], but are notional.

In this Chapter it is introduced a model updating approach in the modal domain for flexiblelink multibody models based on the equivalent rigid link system (ERLS) [74] and formulated through ordinary differential equations (ODEs). The proposed technique starts from the model linearization, in order to apply linear modal analysis, and tackles the peculiarities of the model of FLMSs to ensure a correct comparison between the experimental measurements and the model. In order to provide a more accurate updating, both the measured eigenfrequencies and mode shapes are employed. A numerical technique based on convex optimization is then proposed, to ensure convergence to global optimal solution while accounting for bounds on the feasible values of the updated inertial and elastic parameters.

### 5.2. Model based on the equivalent rigid-link system

### 5.2.1. Definitions

Let us consider an arbitrary FLMS with holonomic and scleronomous constraints. The motion of each link can be notionally separated into the large rigid-body motion of an equivalent rigidlink system (ERLS) and the small elastic displacements of the link with respect to the ERLS itself [74]. Basically, the ERLS represents a moving reference configuration from which the elastic displacements are measured by taking advantage of the FE representations. Such an approach allows for a simple formulation and solution of the nonlinear differential equations governing the system motion. Indeed, once the mutual dynamic coupling between the ERLS and the elastic displacements is correctly accounted for, the ERLS can be studied by taking advantage of the kinematics of rigid body systems as a function of its generalized coordinates q. Additionally, a suitable definition of the ERLS allows to place the moving reference very close to the actual deformed system and hence ensures correctness of the small displacement assumption.

The flexible links of the mechanism are discretized using classical FEs. The absolute position vector $\mathbf{b}_{i}$ of the nodes belonging to the $i$-th FE can be expressed with reference to a global reference frame $\{X, Y, Z\}$ as the sum of two contributions:

$$
\begin{equation*}
\mathbf{b}_{i}=\mathbf{r}_{i}+\mathbf{u}_{i} \tag{5.1}
\end{equation*}
$$

where:

- $\quad \mathbf{u}_{i}$ is the vector of the nodal elastic displacements of the $i$-th element expressed with respect to the ERLS;
- $\quad \mathbf{r}_{i}$ is the vector of the nodal positions of the $i$-th element of the ERLS and it is a function of the ERLS generalized coordinates $\mathbf{q}$, i.e. it can be computed through the rigid-body kinematics:

$$
\begin{equation*}
\mathbf{r}_{i}=\mathbf{r}_{i}(\mathbf{q}) \tag{5.2}
\end{equation*}
$$

The rigid-body kinematics in Eq. (5.2) leads to the following relation between the velocities and the accelerations of the ERLS generalized coordinates and those of the nodes of the ERLS

$$
\begin{align*}
\dot{\mathbf{r}} & =\mathbf{S}(\mathbf{q}) \dot{\mathbf{q}} \\
\ddot{\mathbf{r}} & =\mathbf{S}(\mathbf{q}) \ddot{\mathbf{q}}+\dot{\mathbf{S}}(\mathbf{q}, \dot{\mathbf{q}}) \dot{\mathbf{q}} \tag{5.3}
\end{align*}
$$

where $\mathbf{S}$ is the sensitivity coefficient matrix for all the nodes of the ERLS.


Figure 5.1. ERLS definitions.
Similarly to Eq. (5.1), the position vector of a generic point of the $i$-th finite element $\mathbf{p}_{i}$ is given as the sum of the position of the corresponding point of the ERLS, denoted $\mathbf{w}_{i}$, and the elastic displacements defined with respect to the ERLS, denoted $\mathbf{s}_{i}$ :

$$
\begin{equation*}
\mathbf{p}_{i}=\mathbf{w}_{i}+\mathbf{s}_{i} \tag{5.4}
\end{equation*}
$$

In accordance with the FE theory $\mathbf{s}_{i}$ is computed as follows:

$$
\begin{equation*}
\mathbf{s}_{i}=\mathbf{R}_{i}(\mathbf{q}) \mathbf{N}_{i}\left(\mathrm{x}_{i}, \mathrm{y}_{i}, \mathrm{z}_{i}\right) \mathbf{T}_{i}(\mathbf{q}) \mathbf{u}_{i} \tag{5.5}
\end{equation*}
$$

where

- $\mathbf{N}_{i}$ is the shape function matrix of the FE adopted and it is defined in the local reference frame of the $i$-th element $\left\{x_{i}, y_{i}, z_{i}\right\}$ which follows the motion of the ERLS,
- $\quad \mathbf{R}_{i}$ and $\mathbf{T}_{i}$ are two block-diagonal matrices expressing, respectively, the transformation from the local reference frame to the global one and vice-versa.

A schematic representation of such kinematic definitions is provided in Figure 5.1: a planar case with two-node elements is considered to simplify the drawing.

Let $\mathbf{u}$ be the vector collecting the elastic displacement of all the nodes, and analogously let $\mathbf{b}$ be the vector of the position of all the nodes of the FLMS. Equations (5.1) and (5.3) lead to the following expression that holds in the case of infinitesimal nodal displacements $d \mathbf{b}, d \mathbf{r}$, and $d \mathbf{u}$ :

$$
d \mathbf{b}=d \mathbf{r}+d \mathbf{u}=\left[\begin{array}{ll}
\mathbf{I} & \mathbf{S}
\end{array}\right]\left\{\begin{array}{l}
d \mathbf{u}  \tag{5.6}\\
d \mathbf{q}
\end{array}\right\}
$$

where $\mathbf{I}$ is the identity matrix. The resulting matrix $\left[\begin{array}{ll}\mathbf{I} & \mathbf{S}\end{array}\right]$ is not square since it has $n_{u^{\prime}}$ rows and $n_{u^{\prime}}+n_{q}$ columns, with $n_{u^{\prime}}$ and $n_{q}$ being the numbers of the elastic dofs and of the ERLS generalized coordinates respectively. This means that the linear system in Eq. (5.6) is underdetermined and hence any arbitrary given configuration $d \mathbf{b}$ corresponds to infinitely many sets of $\left\{\begin{array}{ll}\mathbf{u}^{\mathrm{T}} & d \mathbf{q}^{\mathrm{T}}\end{array}\right\}^{\mathrm{T}}$. In order to overcome such an issue, matrix $\left[\begin{array}{ll}\mathbf{I} & \mathbf{S}\end{array}\right]$ should be made square by forcing to zero $n_{q}$ elements of $\mathbf{u}$. In this way, if the resulting matrix is not singular for all the configurations of the mechanism, the system in Eq. (5.6) is determined and hence the position of the ERLS is defined univocally with respect to the actual deformed mechanism.

Let $d \mathbf{u}$ be partitioned into its independent part ( $d \mathbf{u}_{i n} \in \mathbb{R}^{n_{u}=n_{u}-n_{q}}$ ) and into its zeroed part ( $\left.d \mathbf{u}_{0} \in \mathbb{R}^{n_{q}}\right)$ and let $\mathbf{S}$ be accordingly partitioned into $\mathbf{S}_{\text {in }}$ and $\mathbf{S}_{0}$, then Eq. (5.6) can be rewritten as:

$$
d \mathbf{b}=\left[\begin{array}{cc}
\mathbf{I} & \mathbf{S}_{i n}  \tag{5.7}\\
\mathbf{0} & \mathbf{S}_{0}
\end{array}\right]\left\{\begin{array}{c}
d \mathbf{u}_{i n} \\
d \mathbf{q}
\end{array}\right\}
$$

Hereafter, the subscript in will be omitted and symbol $\mathbf{u}$ will have the meaning of independent elastic coordinates.

### 5.2.2. The nonlinear model

The governing equations of motion are obtained by expressing the dynamic equilibrium of the system through the principle of virtual work by splitting the total virtual work into the elemental contributions:

$$
\begin{equation*}
\sum_{i} \int_{v_{i}} \delta \delta_{i}^{\mathrm{T}} \ddot{\mathbf{p}}_{i} \rho_{i} d v+\sum_{i} \int_{v_{i}} \delta \varepsilon_{i}^{\mathrm{T}} \mathbf{D}_{i} \varepsilon_{i} d v=\sum_{i} \int_{v_{i}} \delta \mathbf{p}_{i}^{\mathrm{T}} \mathbf{g} \rho_{i} d v+\left(\delta \mathbf{u}^{\mathrm{T}}+\delta \mathbf{r}^{\mathrm{T}}\right) \mathbf{v} \tag{5.8}
\end{equation*}
$$

where:

- $\boldsymbol{\varepsilon}_{i}$ is the strain vector for the $i$-th element;
- $\quad \mathbf{D}_{i}$ is the stress-strain matrix for the $i$-th element;
- $\quad \rho_{i}$ is the mass density for the $i$-th element;
- g is the gravity acceleration vector;
- v is the vector of the concentrated external forces and torques;
- $\delta \mathbf{u}$ and $\delta \mathbf{r}$ refer to the virtual displacements of all the nodes of the model.

The final expression of the system equations of motion is given by the following set of nonlinear ordinary differential equations (ODEs):

$$
\begin{align*}
& {\left[\begin{array}{cc}
\mathcal{M}_{\mathbf{e}} & \mathcal{M}_{\mathbf{e}} \mathbf{S} \\
\mathbf{S}^{T} \mathcal{M}_{\mathbf{e}} & \mathbf{S}^{T} \mathcal{M}_{\mathbf{e}} \mathbf{S}
\end{array}\right] }  \tag{5.9}\\
&+\left[\begin{array}{cc}
\ddot{\mathbf{u}}(t) \\
\ddot{\mathbf{q}}(t)
\end{array}\right\}+\left[\begin{array}{cc}
\mathbf{2} \mathcal{M}_{\mathbf{e}}+\alpha \mathcal{M}_{\mathrm{e}}+\beta \mathcal{K}_{\mathbf{e}} & -\mathcal{M}_{\mathbf{e}} \dot{\mathbf{S}} \\
\mathbf{0} & \mathbf{0}
\end{array}\right]\left\{\begin{array}{l}
\mathbf{u}(t) \\
\mathbf{S}^{T}\left(\mathbf{2} \mathcal{M}_{G}+\alpha \mathcal{M}_{\mathrm{e}}\right) \\
\mathbf{q}(t)
\end{array}\right\}=\left[\begin{array}{c}
\mathbf{I} \\
\mathbf{S}^{T} \mathcal{M}_{\mathbf{e}} \mathbf{S}
\end{array}\right]\{\mathbf{v}(t)\}
\end{align*}
$$

In Eq. (5.9), matrix $\mathscr{M}_{G}=\mathscr{M}_{G}(\mathbf{q}, \dot{\mathbf{q}})$ represents the centrifugal and Coriolis effects of all the elements, while $\mathcal{M}_{\mathrm{e}}=\mathcal{M}_{\mathrm{e}}(\mathbf{q})$ and $\mathcal{K}_{\mathrm{e}}=\mathcal{K}_{\mathrm{e}}(\mathbf{q})$ are the matrices obtained assembling the consistent mass and stiffness matrices of all the FEs. In Eq. (5.9) damping is modelled by means of the Rayleigh coefficients $\alpha$ and $\beta$.The presence of the off-diagonal terms in the matrices in Eq. (5.9) clearly highlights the coupling between the ERLS generalized coordinates and the elastic displacements. Hence, the dynamic behaviour of the ERLS is not assumed independent from vibration and vice-versa, and the ERLS position is set by the elastic behaviour of the mechanism.

### 5.2.3. The linearized model

Even if the dynamics of flexible link mechanisms and manipulators is nonlinear, it is widely recognized that linearized models are very satisfactory in the case of small deformations about operating points and can be successfully used in the synthesis of control schemes and state observers [79]. On the other hand, linear models represented through ODEs allow for the straightforward application of the modal analysis. Hence, updating the model of a FLMS by means of linearized model is an effective mean for the identifications of the parameters of the non-linear model, which is instead not suitable for the modal analysis.

As for the use of ODEs for the model synthesis, despite the large use of differential algebraic equations (DAEs) for modelling FLMSs, it is again justified by the possibility to employ the modal analysis. Indeed, models employing DAEs, which are formulated through second-order ordinary differential equations coupled to a set of algebraic constraint equations, do not allow for such a kind of analysis. More complicate and computational challenging methods should be therefore adopted to handle this kind of mathematical formulation [80].

In order to linearize the model, let us consider small displacements about a static equilibrium configuration, which is set by the equilibrium configuration of the ERLS $\mathbf{q}_{\mathrm{e}}$. Then the following linear model is obtained by just considering the first-order terms of the Taylor's expansion:

$$
\begin{align*}
& +\left[\begin{array}{cc}
\mathcal{K}_{e} & \mathbf{0} \\
\mathbf{0} & -\left(\frac{\mathrm{d}\left(\mathbf{S}^{T} \mathcal{M}_{\mathrm{e}}\right)}{\mathrm{d} \mathbf{q}} \otimes \mathbf{g}+\frac{\mathrm{d} \mathbf{S}^{T}}{\mathrm{dq}} \otimes \mathbf{v}_{\mathbf{e}}\right)
\end{array}\right]_{\mathbf{q}=\mathbf{q}_{e}}\left\{\begin{array}{c}
\mathbf{u}(t) \\
\mathbf{q}(t)
\end{array}\right\}=\left[\begin{array}{c}
\mathbf{I} \\
\mathbf{S}^{T}
\end{array}\right]_{\mathbf{q}=\mathrm{q}_{e}}\{\mathbf{v}\} \tag{5.10}
\end{align*}
$$

In Eq. (5.10), the matrices $\mathcal{M}_{\mathrm{e}}, \mathbf{S}, \mathcal{K}_{\mathrm{e}}$ and the derivatives $\mathrm{d}\left(\mathbf{S}^{T} \mathcal{M}_{\mathrm{e}}\right) / \mathrm{d} \boldsymbol{q}$ and $\mathrm{d} \mathbf{S}^{T} / \mathrm{d} \boldsymbol{q}$ are computed about the equilibrium configuration. The term $\left(\mathrm{d} \mathbf{S}^{\mathrm{T}} / \mathrm{dq}\right) \otimes \mathbf{v}_{\mathrm{e}}$ represents the inner product of matrix $\left[\partial S_{i, 1} / \partial q_{j} \cdots \partial S_{i, n} / \partial q_{j}\right]$ with vector $\mathbf{v}_{\mathbf{e}}$, for all the subscripts $i, j$ (obvious extension to $\left.\left(\mathrm{d}\left(\mathbf{S}^{T} \mathcal{M}_{\mathrm{e}}\right) / \mathrm{d} \mathbf{q}\right) \otimes \mathbf{g}\right)$.

Since model updating should be performed by investigating the system in a stable equilibrium configuration, gravity forces should be compensated in the experimental analysis whenever the mechanism lies in the vertical plane. Indeed, the absence of a stable equilibrium would make the mechanism diverge from the initial configuration after the excitation and therefore the system dynamics would not meet the hypothesis of the linearized model. One of the possible approaches is to balance gravity force through tuned external springs that determine asymptotic stability of the equilibrium point. Clearly, springs should be then included in the stiffness matrix of the model. In contrast the use of brakes, such as those of the actuators, modifies the boundary conditions and makes the mechanism behave as a structure by removing all the model terms representing the coupling between the ERLS and the elastic displacements. Therefore model updating may lead to less reliable results.

### 5.3. Data consistency between model and experimental modal analysis

### 5.3.1. Size compatibility between the FE model and the measurements

Model updating requires one-to-one correspondence between the model coordinates and measurement data. From a practical view point this requirement is usually not satisfied because of the reduced number of sensors used in the experimental measurements and the presence of some inaccessible locations. Additionally, rotational dofs cannot be easily measured. Thus, the displacement vector of the experimental data is usually not compatible with the FE model of the FLMS and hence, as a preliminary step before the model updating, it should be properly modified.

The most suitable approach, that is the one adopted in this work, is expanding the experimental data to match the dimension of the FE model [76, 81, 82]. The use of curve fitting, by means of the shape functions of the FEs, is for example an effective and simple way.

### 5.3.2. Data representation from the physical reference frame to the ERLS reference frame

The second relevant issue that should tackled is the transformation of the measured eigenvectors from the reference adopted to measure the mode shapes, which is physical, to the fictitious one
adopted in the ERLS-based FE model. Indeed, the sensors employed to perform experimental modal analysis tests, such as for example the accelerometers, typically provide measurements with respect to local reference frames corresponding to the mechanism in its initial static configuration. Conversely the displacement vector of the FLMS model (u) includes the elastic displacements with respect to the ERLS and expressed in a global reference frame. Hence, in accordance with Eq. (5.7) and with the final model in Eq. (5.10), some of the elastic displacements are forced to zero to make the ERLS follow the actual elastic mechanism. At the same time, the ERLS coordinates should be included in the displacement vector and therefore in the eigenvectors.

Let $\mathbf{q}_{e}$ be the static equilibrium configuration corresponding to the actual position of the mechanism, adopted for both the experimental modal analysis and the model linearization. Let us assume that the mechanism is vibrating in accordance with the $i$-th mode about the equilibrium configuration. The absolute position of the nodes of the FEM, denoted $\mathbf{b}_{i}$, can be represented both through the eigenvector defined in the physical frame adopted for the experimental analysis, $\overline{\boldsymbol{\eta}}_{i}$, and through the one defined with the model coordinates, $\boldsymbol{\eta}_{i}$. If $\mathbf{b}_{i}$ is defined in a global frame, the eigenvector measured in the physical frame should be rotated onto the global reference frame through the ERLS kinematic constraint equations:

$$
\begin{equation*}
\mathbf{b}_{i}=\overline{\mathbf{r}}\left(\mathbf{q}_{e}\right)+\boldsymbol{\Theta}\left(\mathbf{q}_{e}\right) \overline{\boldsymbol{\eta}}_{i} \tag{5.11}
\end{equation*}
$$

$\overline{\mathbf{r}}$ is the vector of the nodal position of the FLMS in the configuration set by $\mathbf{q}_{e}$, while $\boldsymbol{\Theta}$ is a block transformation matrix which depends on the rigid-body kinematic constraints.

If, in contrast, the same configuration is represented through the eigenvector of the linearized model $\boldsymbol{\eta}_{i}$ expressed with respect to the ERLS at configuration $\mathbf{q}_{i}=\mathbf{q}_{e}+\Delta \mathbf{q}$, the following relation holds

$$
\begin{equation*}
\mathbf{b}_{i}=\mathbf{r}\left(\mathbf{q}_{i}\right)+\boldsymbol{\eta}_{i} \tag{5.12}
\end{equation*}
$$

where $\mathbf{r}\left(\mathbf{q}_{i}\right)$ is the vector of the nodal position of the ERLS in the configuration set by $\mathbf{q}_{i}$. The displacement of the ERLS generalized coordinates $\Delta \mathbf{q}$ represents the rotation of the ERLS
which moves to follow the actual mechanism. Hence it is equal to the $n_{q}$ zeroed elements of $\mathbf{u}$ already discussed in Section 2.1.

A graphical representation of Eq. (5.11) and Eq. (5.12) is provided in Figure 5.2 with reference to a single flexible link.

The final transformation to make the experimental eigenvectors consistent with the model formulation can be immediately inferred from Eqs. (5.11) and (5.12):

$$
\begin{equation*}
\hat{\boldsymbol{\eta}}_{i}=\overline{\mathbf{r}}\left(\mathbf{q}_{e}\right)-\mathbf{r}\left(\mathbf{q}_{i}\right)+\boldsymbol{\Theta}\left(\mathbf{q}_{e}\right) \overline{\boldsymbol{\eta}}_{i} \tag{5.13}
\end{equation*}
$$

where $\hat{\boldsymbol{\eta}}_{i}$ denotes the transformed experimental eigenvector, that is suitable for being employed in the model updating procedure.


Figure 5.2. Interpretation of the analytical and experimental eigenvectors.

### 5.4. The model updating technique

### 5.4.1. Formulation of the mass and stiffness matrix updating

The model is updated by correcting the inertial and elastic parameters through the undamped model. This simplifies the formulation of the inverse eigenvalue problem by casting it as a first order problem and by using real eigenpairs, and then improves its numerical solution by just
handling real numbers. The use of the undamped models is widely adopted in structural dynamics and is justified by the negligible internal damping of the beams and of the low friction of the joints. Hence, the identification of the Rayleigh coefficients can be performed separately, through the well established methods [83, 84], once the inertial an elastic parameters have been updated.

The linearized undamped model in Eq. (5.10) is represented through the mass $\mathbf{M}^{O} \in \mathbb{R}^{N \times N}$ and the stiffness $\mathbf{K}^{o} \in \mathbb{R}^{N \times N}$ matrices, which are the initial (nominal) system matrices based on the nominal inertial and elastic parameters ( $N$ is number of model dofs, including both rigid and elastic coordinates: $N=n_{q}+n_{u}$ ). The updating of the model matrices is here represented through the additive correction matrices, the so called updating matrices $\boldsymbol{\Delta M} \in \mathbb{R}^{N \times N}$ and $\Delta \mathbf{K} \in \mathbb{R}^{N \times N}$ which collect the parameters that are affected by uncertainty and hence should be corrected. The topology of the updating matrices depends on the topology of $\mathbf{M}^{O}$ and $\mathbf{K}^{O}$. The following representations are therefore proposed for $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$ :

$$
\begin{equation*}
\boldsymbol{\Delta} \mathbf{M}=\sum_{j=1}^{N_{m}} \boldsymbol{\Delta} \mathbf{M}_{j}=\sum_{j=1}^{N_{m}}\left[\frac{\partial \mathbf{M}^{O}}{\partial m_{j}^{O}}\right] \Delta m_{j} ; \quad \boldsymbol{\Delta} \mathbf{K}=\sum_{h=1}^{N_{k}} \boldsymbol{\Delta} \mathbf{K}_{h}=\sum_{h=1}^{N_{m}}\left[\frac{\partial \mathbf{K}^{O}}{\partial k_{h}^{O}}\right] \Delta k_{h} \tag{5.14}
\end{equation*}
$$

$\Delta \mathbf{M}_{j}$ and $\Delta \mathbf{K}_{h}$ are the updating sub-matrices related to each updating parameter $m_{j}^{o}$ or $k_{h}^{o}$ that should be updated through the additive corrections $\Delta m_{j}$ and $\Delta k_{h}$.

$$
\begin{equation*}
m_{j}=m_{j}^{O}+\Delta m_{j} ; \quad k_{h}=h_{h}^{O}+\Delta k_{h} \tag{5.15}
\end{equation*}
$$

Hence, the mass and stiffness matrices of the $N$-dimensional updated model are $\mathbf{M}^{o}+\boldsymbol{\Delta} \mathbf{M}$ and $\mathbf{K}^{o}+\Delta \mathbf{K}$. Finally, $N_{m}$ and $N_{k}$ denote the number of parameters that can be updated in the mass and stiffness matrix, respectively.

In order to improve the numerical conditioning of the problem by tackling the different scaling of the inertial and elastic parameters, dimensionless corrections $\bar{m}_{j}$ and $\bar{k}_{h}$ are defined and hereafter adopted in the updating:

$$
\begin{equation*}
\bar{m}_{j}=\frac{\Delta m_{j}}{m_{j}^{O}} ; \quad \bar{k}_{h}=\frac{\Delta k_{h}}{k_{h}^{O}} \tag{5.16}
\end{equation*}
$$

For those $m_{j}^{O}$ or $k_{h}^{O}$ whose value is zero in the nominal model, the reference value might be defined as the upper bound of their allowed modifications. Nonetheless other arbitrary, though reasonable, reference numbers can be assumed. All the variables $\bar{m}_{j}$ and $\bar{k}_{h}$ are collected in the $\left(N_{m}+N_{k}\right)$-dimensional vector of the updating parameters $\chi=\left\{\begin{array}{lllllll}\bar{m}_{1} & \ldots & \bar{m}_{N_{m}} & \bar{k}_{1} & \ldots & \bar{k}_{N_{k}}\end{array}\right\}^{\mathrm{T}}$, which collects all the problem unknown.

By taking advantage of the definitions in Eqs. (5.14), (5.15) and (5.16), the updating matrices can be written as follows:

$$
\begin{equation*}
\boldsymbol{\Delta} \mathbf{M}=\sum_{j=1}^{N_{m}}\left[m_{j}^{O} \frac{\partial \mathbf{M}^{O}}{\partial m_{j}^{O}}\right] \bar{m}_{j}:=\sum_{j=1}^{N_{m}} \mathbf{A}_{j} \bar{m}_{j} ; \quad \mathbf{\Delta} \mathbf{K}=\sum_{h=1}^{N_{k}}\left[k_{h}^{O} \frac{\partial \mathbf{K}^{O}}{\partial k_{h}^{O}}\right] \bar{k}_{h}:=\sum_{h=1}^{N_{k}} \mathbf{B}_{h} \bar{k}_{h} \tag{5.17}
\end{equation*}
$$

where the sub-matrices $\mathbf{A}_{j}$ and $\mathbf{B}_{h}$ have been introduced for brevity of notation. The partial derivatives in Eq. (5.17) are constant whenever Young's modulus, mass density, nodal masses, nodal inertias, sectional area or lumped springs should be updated. In contrast, if some $\bar{m}$ or $\bar{k}$ represent the length of the beams, the derivatives are not constant and therefore iterations should be made in the updating procedure. However, it is often reasonable assuming the beam length as a known and exact parameter.

### 5.4.2. Constraints on the updating parameters

The correction of the model parameters should be also bounded to ensure physical meaning of the numerical solution and to avoid unreasonable solutions. Hence, lower and upper bounds on the single variables (denoted respectively $\chi_{\mathrm{L}}$ and $\chi_{\mathrm{U}}$ ) and on combinations of more variables (such as on the total mass) should be defined. The latter type of constraint is defined through a convex polyhedron (by means of a matrix $\mathbf{G}$ and a vector $\boldsymbol{\gamma}$ ). These two types of constraint define the following convex feasible set $\Gamma$ (where the inequalities are element-wise):

$$
\begin{equation*}
\Gamma=\left\{\chi \quad \mid \quad\left(\chi_{\mathrm{L}} \leq \chi \leq \chi_{\mathrm{U}}\right) \cap \mathbf{G} \chi \leq \gamma\right\} \tag{5.18}
\end{equation*}
$$

### 5.4.3. Formulation of the model updating problem

On the basis of the model assumed for representing the updating matrices, the following eigenvalue problem must hold for any measured eigenpairs $\left(\hat{\omega}_{i}^{2}, \hat{\boldsymbol{\eta}}_{i}\right), i=1, \ldots, n$, where $n \leq N$ is the number of measured vibrational modes and $\hat{\omega}_{i}$ is the natural frequency associated to the eigenvector $\hat{\boldsymbol{\eta}}_{i}$ :

$$
\begin{equation*}
\hat{\omega}_{i}^{2}\left[\mathbf{M}^{o}+\Delta \mathbf{M}\right] \hat{\boldsymbol{\eta}}_{i}-\left[\mathbf{K}^{o}+\Delta \mathbf{K}\right] \hat{\boldsymbol{\eta}}_{i}=\mathbf{0} \tag{5.19}
\end{equation*}
$$

Any arbitrary normalization of the measured eigenvectors can be adopted.
By introducing the representation of the updating matrices given in Eq. (5.17) and premultiplying by $\hat{\boldsymbol{\eta}}_{i}{ }^{\mathrm{T}}$, the eigenvalue problem related to each eigenpair can be written as a scalar equation, which is a linear problem in the unknown $\left\{\bar{m}_{j}, \bar{k}_{h}\right\}$ :

$$
\begin{equation*}
\hat{\boldsymbol{\eta}}_{i}^{T}\left[\hat{\omega}_{i}^{2} \mathbf{M}^{O}-\mathbf{K}^{o}\right] \hat{\boldsymbol{\eta}}_{i}+\hat{\omega}_{i}^{2} \sum_{j=1}^{N_{m}} \boldsymbol{\eta}_{i}^{T} \mathbf{A}_{j} \hat{\boldsymbol{\eta}}_{i} \bar{m}_{j}-\sum_{h=1}^{N_{k}} \boldsymbol{\eta}_{i}^{T} \mathbf{B}_{h} \hat{\boldsymbol{\eta}}_{i} \bar{k}_{h}=0 \tag{5.20}
\end{equation*}
$$

Equation (5.20) can be represented in the typical compact form of a linear algebraic system through the vector $\mathbf{a}_{i}$, which collects the known coefficients, and the scalar $b_{i}=-\hat{\boldsymbol{\eta}}_{i}^{T}\left[\hat{\omega}_{i}^{2} \mathbf{M}^{o}-\mathbf{K}^{o}\right] \hat{\boldsymbol{\eta}}_{i}$ which is the known term:

$$
\begin{equation*}
\mathbf{a}_{i} \chi-b_{i}=0 \tag{5.21}
\end{equation*}
$$

Since the problem unknowns are constrained by the feasible set $\Gamma$ and the condition in Eq. (5.21) should be simultaneously satisfied for all the measured eigenpairs ( $i=1, \ldots, n$ ), an exact solution might not exist. Hence the problem should be approximated as a 2-norm minimization problem (least square minimization):

$$
\begin{equation*}
\min \|\mathbf{W}(\mathbf{A} \boldsymbol{\chi}-\mathbf{b})\|_{2}^{2} \tag{5.22}
\end{equation*}
$$

$\mathbf{A} \in \mathbb{R}^{n \times\left(N_{m}+N_{k}\right)}$ and $\mathbf{b} \in \mathbb{R}^{n}$ collect the terms of each eigenvalue problem in Eq. (5.21): the rows of matrix $\mathbf{A}$ are the vectors $\mathbf{a}_{i}$ while the entries of $\mathbf{b}$ are the scalar coefficients $b_{i}$. The diagonal matrix $\mathbf{W} \in \mathbb{R}^{n \times n}$ is defined to weight the residual of each eigenvalue problem through the
positive and scalar weights $w_{i}$, with $i=1, \ldots, n$. The introduction of the weighing matrix has several practical advantages. Indeed it allows to give less importance to those eigenpairs whose measurements are less reliable, e.g. because of the presence of noise or uncertainty, as well as to give higher importance to the vibrational modes that are more observable and controllable and hence provide more relevant contributions to the forced response.

The numerical solution of the minimization problem should be improved through preconditioning to compensate for the simultaneous presence of low-frequency and highfrequency modes. The use of the pseudo-inverse matrix (denoted with the superscript ${ }^{\dagger}$ ) as the preconditioner allows to write the system in Eq. (5.22) in the following equivalent form, whose condition number is one:

$$
\begin{equation*}
\min \left\|\chi-[\mathbf{W} \mathbf{A}]^{\dagger} \mathbf{W} \mathbf{b}\right\|_{2}^{2} \tag{5.23}
\end{equation*}
$$

Finally, regularization can be adopted to penalize selectively the components of $\chi$ [85]. Indeed, some model parameters might be more uncertain than the others and therefore it is more desirable the updating procedure to modify them. Therefore, the term $\lambda\|\boldsymbol{\Omega} \boldsymbol{\chi}\|_{2}^{2}$ is added to the minimization problem in Eq. (5.23). The scalar value $\lambda$ is named the regularization parameter and trades-off between the cost of missing the eigenpair specifications (i.e. Eq. (5.21)) and the cost of using large values of the design variables (i.e. $\|\chi\|_{2}^{2}$ ). The positive-definite real matrix $\boldsymbol{\Omega} \in \mathbb{R}^{\left(N_{m}+N_{k}\right) \times\left(N_{m}+N_{k}\right)}$ is the regularization operator and can be adopted to define the relative weight between the updates of the different model parameters.

The model updating problem is therefore finally represented as the following constrained minimization problem:

$$
\left\{\begin{array}{c}
\min _{\chi}\left\|\left[\begin{array}{c}
\mathbf{I} \\
\lambda \boldsymbol{\Omega}
\end{array}\right] \chi-\left\{\begin{array}{c}
{[\mathbf{W A}]^{\dagger} \mathbf{W b}} \\
\mathbf{0}
\end{array}\right\}\right\|_{2}^{2}  \tag{5.24}\\
\chi \in \Gamma
\end{array}\right.
$$

Both the objective function and the feasible domain are convex. Hence model updating is cast as a convex, constrained minimization problem whose solution is straightforward. Indeed the
problem convexity ensures that the global optimal solution can be found regardless of the selection of the initial guess required by the numerical minimization. In addition, several reliable and efficient numerical algorithms are available in commercial software.

### 5.5. Experimental validation

### 5.5.1. Description of the experimental setup

The model updating strategy proposed has been applied to the manipulator with flexible links shown in Figure 5.3(a). Such a manipulator consists of five rods with circular cross-sections made of Aluminum Anticorodal. The links are connected to each other by means of revolute joints, created by means of double ball bearings with negligible clearance and friction losses, inserted in Aluminum housings. A hybrid serial-parallel kinematic topology is obtained. Links 1,2 and 5 are driven by three brushless motors. The mechanism is equipped with two balancing springs, modelled as a single lumped spring, which are employed to reduce the static torque required by the servomotors to compensate for gravity force.

All the links have been modelled with two-node and six-degree-of-freedom Euler-Bernoulli's beam elements to represent the planar motion of the mechanism, overall nine beam elements have been adopted (see Figure 5.3(b)). Lumped masses and nodal inertias have been adopted to represent the joints and the motors. The resulting model has 27 elastic dofs and three ERLS generalized coordinates $\mathbf{q}=\left\{\begin{array}{lll}q_{1} & q_{2} & q_{3}\end{array}\right\}^{T}$, which are the absolute rotations of the actuators. The rigid moving reference is then defined by forcing the ERLS to coincide with the actual flexible mechanism at the rotation of these three joints. Hence, three elastic displacements with respect to the ERLS are forced to zero in the FLMS model, in order to place correctly the moving reference defined by the equivalent rigid mechanism (see Section 2 for the theoretical discussion of this issue).

Model updating has been done by investigating the mechanism in the stable equilibrium configuration, where gravity forces are exactly compensated by the balancing force exerted by the two external springs, which corresponds to $\mathbf{q}_{e}=\left\{\begin{array}{lll}1.0821 & 2.6354 & 4.7123\end{array}\right\}^{\mathrm{T}}[\mathrm{rad}]$.

Experimental modal analysis has been performed through impact tests. A series of frequency response functions (FRF) have been measured at various positions through accelerometers, while an instrumented impact hammer has been adopted to supply the input forces.

The measurement chain consists in:

- 5 accelerometers, whose specifications are listed in Table 5.1;
- an instrumented impact hammer PCB Modally Tuned $086 C 03$ with teflon tip;
- a frontend system LMS SCADAS Recorder (max input $\pm 10 \mathrm{~V}$, ADC 24 bit);
- a modal software LMS Test.Lab Impact Testing 11B, with tools \{Geometry, Modal Analysis, PolyMAX Modal Analysis $\}$.

The impact setup defined for the test is described in Table 5.2.
The number of measurement points has been selected to allow for a reliable representation of the mode shapes in the frequency range of interest and for the interpolation to expand the experimental results while smoothing measurement noise and uncertainty. In accordance with these aims, a manipulator geometry has been created through 22 points within the software tool LMS Test.Lab Impact Testing 11B Geometry, as shown in Figure 5.4.


Figure 5.3. (a) picture of the planar, three-dof manipulator and (b) its finite element model.

Table 5.1. Accelerometers employed.

| Mandfacturer | Model | Typology | Sensitivity [mV/g ] |
| :---: | :---: | :---: | :---: |
| Endevco | 27AM1-100 10203 |  |  |
| Dytran | $3136 A 1945$ | ICP | 102.4 |
| PCB | ICP | 100 |  |
| Brüel \& Kjær | 4508 | Full bridge DC | 66.7 |
| Brüel \& Kjær | 4506 B | ICP | 9.78 |
|  |  | X-triaxial ICP | 93.5 |

Table 5.2. Impact test settings.

| ACQUISITION SETTING |  |
| :--- | :---: |
| Sample frequency | 2048 Hz |
| Frequency resolution | 0.06 Hz |
| Acquisition time | 16 s |
| TriGGERING |  |
| Trigger level | 5.72 N |
| Pretrigger | 0.01 s |
| WINDOWING |  |
| Input | Uniform |
| Response | Exponential (decay 90\%) |
| MEASUREMENT AVERAGES | 3 |



Figure 5.4. Manipulator geometry created by the GUI of LMS Test.Lab 11B.

As far as the frequency range of interest is concerned, it ranges from 0 to 200 Hz . This interval corresponds to the bandwidth of the motor current loop and hence to the range of frequencies
where active vibration control could be reasonably achieved. In contrast, the vibrational modes with higher frequency cannot be neither controlled nor excited in the forced motion, and hence they are not of interest for the model. Therefore, seven vibrational modes have been adopted to perform model updating.

A total number of 315 FRFs has been recorded by averaging three measurements. All the acquisitions are summarized in the experimental frequency response matrix shown in Figure 5.5. The $X$ and $Y$ directions indicated in Figure 5.5 are referred to the local reference systems of each link, as plotted in Figure 5.4. As a representative example of such FRFs, Figure 5.6 shows the 7th row of the matrix in Figure 5.5. Data have been post processed using the tool PolyMAX of LMS Test.Lab 11B, which extracts the mode shapes exploiting the so-called poleresidual model [86].

|  | $$ | $\begin{aligned} & \stackrel{N}{\square} \\ & \underset{\square}{B} \end{aligned}$ | $\begin{aligned} & \frac{?}{\square} \\ & \frac{n}{B} \end{aligned}$ |  | $\begin{aligned} & \overline{\mathrm{x}} \\ & \underset{\Xi}{\Xi} \end{aligned}$ | $\begin{aligned} & \text { n } \\ & \underset{\sim}{y} \\ & \vdots \end{aligned}$ | $\begin{aligned} & \stackrel{~}{\mathrm{y}} \\ & \underset{Z}{B} \end{aligned}$ | $\begin{aligned} & \overline{\ddot{n}} \\ & \vec{B} \end{aligned}$ | $\begin{aligned} & \tilde{y} \\ & \dot{\tilde{n}} \end{aligned}$ | $\begin{aligned} & \ddot{\tilde{n}} \\ & \underset{B}{B} \end{aligned}$ | $\frac{\vec{ت}}{\vec{j}}$ | $\begin{aligned} & \stackrel{y}{\dot{G}} \\ & : \neq \end{aligned}$ | $\begin{aligned} & \stackrel{?}{\dot{G}} \\ & \frac{\square}{\Xi} \end{aligned}$ | $\begin{aligned} & n \\ & \stackrel{n}{甘} \\ & \vdots \end{aligned}$ | $\begin{aligned} & \stackrel{0}{\dot{G}} \\ & \dot{B} \end{aligned}$ | $\begin{aligned} & \underset{\dot{G}}{\underset{B}{B}} \\ & \hline \end{aligned}$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| OUTPUT | Dir | Y | Y | Y | Y | Y | Y | Y | Y | Y | Y | Y | Y | Y | Y | Y |
| link 1:2 | Y | - - + | - - + | -1+ | -\+ | -1+ | - $1+$ | - - + | - $1+$ | - - + | - - | - - | - - | - - | - - | - - |
| link 1:3 | Y | + + | + + | + - | + + | $+1+$ | $+\backslash$ | $+1+$ | + + | $+1+$ | $+1+$ | + + | + + | + - | + - | +\+ |
| link1:4 | Y | $+1+$ | +- | + + | $+1+$ | +1+ | $+1+$ | + + | + + | $+1+$ | $+\backslash+$ | +\- | + + | +- | +- | $+\+$ |
| link 1:5 | Y | -1+ | - $1+$ | -1+ | - + + | -1+ | - $1+$ | -1- | - $1+$ | -1+ | - - + | -1+ | - $1+$ | -1- | - - | - - |
| link2:1 | Y | +1+ | + + | + + | $+1+$ | + + | $+1+$ | +\- | $+1+$ | + + | + + | +1+ | $+1+$ | +- | + - | + - |
| link2:2 | Y | $+1+$ | + + | + - | $+\backslash+$ | $+1+$ | $+1+$ | $+1+$ | + + | $+1+$ | $+1+$ | + + | + + | + - | + - | + + + |
| link2:3 | Y | $+1+$ | +\- | + + | $+1+$ | +1+ | $+1+$ | +1+ | + + | +1+ | $+1+$ | +\- | $+1+$ | + | +- | $+1+$ |
| link2:4 | Y | $+1+$ | $+1+$ | + + + | $+\backslash+$ | $+1+$ | + + | + + | + + | $+1+$ | + - | +\- | + - | + - | $+-$ | + - |
| link3:1 | Y | -1+ | -l+ | -1+ | - - + | - - + | -l+ | -1- | -1+ | - - + | -l+ | -1+ | - - + | -1- | -- | -- |
| link3:2 | Y | - $1+$ | - - | -1+ | -\+ | -1+ | - $1+$ | -1+ | - $1+$ | - - + | - $1+$ | -- | - $1+$ | -1- | - - | - - + |
| link3:3 | Y | - + | - $1+$ | -- | - - + | - + + | -l+ | - - + | - ++ | - - + | - - + | - - + | -l+ | -- | -- | - - + |
| link3:4 | X | +\- | +1- | +1- | + + | + + | + + | +1+ | + + | + + | + + + | +1+ | + + | +1+ | + + + | + + + |
| link4:1 | Y | - - + | - - + | -1+ | - - + | - - + | - $1+$ | - - + | - $1+$ | - - + | - - | -1- | - - | -1- | - - | -l- |
| link4:2 | Y | -1+ | - - | -1+ | -1+ | -1+ | - $1+$ | -1+ | - $1+$ | - - + | - - + | -1- | -1+ | -1- | - - | - - + |
| link4:3 | Y | - + + | -\+ | -1+ | -\+ | -l+ | - + | - - + | -l+ | - - + | - - | -- | -- | -1- | -- | -l- |
| link4:4 | Y | -1+ | -l+ | -- | -l+ | -l+ | -l+ | - - + | -l+ | -l+ | - - + | -1+ | -l+ | -1- | -- | -I+ |
| link4:5 | Y | - $1+$ | - - + | -1+ | - - + | - - + | -l+ | - - | - $1+$ | - - + | - + + | -1+ | - $1+$ | -1- | - - | -- |
| link4:6 | Y | + - | +- | +- | + + + | + + | + + | +1+ | + + | + + | + + | + + | + + | +1+ | + + | + + |
| link4:7 | X | -1+ | - - + | -1+ | - - + | -1+ | - $1+$ | -1- | - $1+$ | -1+ | -l+ | -1+ | - $1+$ | -1- | - - | -l- |
| link4:7 | Y | - - + | -l+ | -1+ | -l+ | - - + | -l+ | -1- | -1+ | - - + | -1+ | - - + | -1+ | -1- | -- | -- |
| link5:2 | Y | - - | - - | -- | - $1+$ | - - + | - $1+$ | - - + | - $1+$ | - - + | - - + | -1+ | - $1+$ | -1+ | - - + | $-1+$ |

Figure 5.5. Frequency response matrix.


Figure 5.6. 7th row of the frequency response matrix.

### 5.5.2. Size compatibility between the FE model and the measurements

Coordinate expansion has been adopted to ensure size compatibility between the model and the measurements. In particular, least-square curve fitting has been employed to estimate the unmeasured rotational dofs, while filtering measurement noise. Fitting is based on the measurement of a redundant set of measurable translational dofs and on the least-square regression of the measured displacements along the links by means of the polynomial interpolation shape functions of the Euler-Bernoulli's beam elements.

### 5.5.3. Statement of the model updating problem

The updating parameters collected in $\chi$ have been defined as the lumped stiffness of the balancing springs ( $\mathrm{K}_{\mathrm{bs}}$ ), the nodal inertias (J) and masses (M), the mass density ( $\rho$ ) of the links (which is assumed equal among all the links), the Young's modulus (E) (which has been instead treated separately for each link to compensate for local stiffening due to the kinematic joints).The lengths of the links and their cross sectional area are instead assumed as exactly
known. The feasible values have been constrained by upper and lower bounds, which are listed in Table 5.3 together with the original and the updated values for each parameter. Additionally, a constraint on the total mass has been included, by formulating it through Eq. (5.18). The following notation has been employed for the subscripts in Table 5.3: numbers from 1 to 5 denote the links (see Figure 5.3(b)), letters from A to E indicate the joints, H indicates the lumped mass of the balancing spring while $F$ the one attached to link 5; the number adjacent to the letter denotes the link which each lumped mass is connected to. A total number of 29 parameters has been updated.

Table 5.3. Updating parameters.

| INDEX | PARAMETER | NOMINAL VALUE | LOWER BOUND | UPPER BOUND | UPDATE <br> VALUE |
| :---: | :---: | :---: | :---: | :---: | :---: |
| 1 | $\rho\left[\mathrm{kg} / \mathrm{m}^{3}\right]$ | 2710 | 2574.5 | 2845.5 | 2742.6 |
| 2 | $\mathrm{E}_{1}[\mathrm{~Pa}]$ | $69 \cdot 10^{9}$ | $65.55 \cdot 10^{9}$ | $72.45 \cdot 10^{9}$ | $65.55 \cdot 10^{9}$ |
| 3 | $\mathrm{E}_{2}[\mathrm{~Pa}]$ | $69 \cdot 10^{9}$ | $65.55 \cdot 10^{9}$ | $72.45 \cdot 10^{9}$ | $65.55 \cdot 10^{9}$ |
| 4 | $\mathrm{E}_{3}[\mathrm{~Pa}]$ | $69 \cdot 10^{9}$ | $71.07 \cdot 10^{9}$ | $72.45 \cdot 10^{9}$ | $71.07 \cdot 10^{9}$ |
| 5 | $\mathrm{E}_{4}[\mathrm{~Pa}]$ | $69 \cdot 10^{9}$ | $69.69 \cdot 10^{9}$ | $72.45 \cdot 10^{9}$ | $69.69 \cdot 10^{9}$ |
| 6 | $\mathrm{E}_{5}[\mathrm{~Pa}]$ | $69 \cdot 10^{9}$ | $65.55 \cdot 10^{9}$ | $72.45 \cdot 10^{9}$ | $68.99 \cdot 10^{9}$ |
| 7 | $\mathrm{K}_{\mathrm{bs}}[\mathrm{N} / \mathrm{m}]$ | 1520 | 1216 | 1824 | 1520 |
| 8 | $\mathrm{M}_{\mathrm{A}, 1}[\mathrm{~kg}]$ | 7.6440 | 6.1152 | 9.1728 | 7.6509 |
| 9 | $\mathrm{M}_{\mathrm{B}, 1}[\mathrm{~kg}]$ | 0.3920 | 0.3136 | 0.4704 | 0.3927 |
| 10 | $\mathrm{M}_{\mathrm{A}, 2}[\mathrm{~kg}]$ | 9.5170 | 7.6136 | 11.4204 | 9.5190 |
| 11 | $\mathrm{M}_{\mathrm{H}, 2}[\mathrm{~kg}]$ | 0.4000 | 0.3200 | 0.4800 | 0.4000 |
| 12 | $\mathrm{M}_{\mathrm{D}, 2}[\mathrm{~kg}]$ | 0.6480 | 0.5184 | 0.7776 | 0.6484 |
| 13 | $\mathrm{M}_{\mathrm{B}, 3}[\mathrm{~kg}]$ | 0.6710 | 0.5368 | 0.8052 | 0.6732 |
| 14 | $\mathrm{M}_{\mathrm{C}, 3}[\mathrm{~kg}]$ | 0.3830 | 0.30644 | 0.4596 | 0.3837 |
| 15 | $\mathrm{M}_{\mathrm{C}, 4}[\mathrm{~kg}]$ | 0.6580 | 0.5264 | 0.7896 | 0.6601 |
| 16 | $\mathrm{M}_{\mathrm{D}, 4}[\mathrm{~kg}]$ | 0.3080 | 0.2464 | 0.3696 | 0.3080 |
| 17 | $\mathrm{M}_{\mathrm{E}, 4}[\mathrm{~kg}]$ | 1.5370 | 1.8444 | 2.1518 | 1.8444 |
| 18 | $\mathrm{M}_{\mathrm{E}, 5}[\mathrm{~kg}]$ | 0.0950 | 0.1140 | 0.1330 | 0.1140 |
| 19 | $\mathrm{M}_{\mathrm{F}, 5}[\mathrm{~kg}]$ | 0.0460 | 0.0368 | 0.0552 | 0.0460 |
| 20 | $\mathrm{J}_{\mathrm{A}, 1}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0.0130 | 0.0104 | 0.0156 | 0.01302 |
| 21 | $\mathrm{J}_{\mathrm{A}, 2}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0.0229 | 0.0183 | 0.0275 | 0.0229 |
| 22 | $\mathrm{J}_{\mathrm{E}, 5}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0.0015 | 0.0012 | 0.0018 | 0.0015 |
| 23 | $\mathrm{J}_{\mathrm{B}, 1}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0 | 0 | $4 \cdot 10^{-4}$ | $4 \cdot 10^{-4}$ |
| 24 | $J_{\text {B, } 3}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0 | 0 | $4 \cdot 10^{-4}$ | $4 \cdot 10^{-4}$ |
| 25 | $\mathrm{J}_{\mathrm{C}, 3}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0 | 0 | $4 \cdot 10^{-4}$ | $4 \cdot 10^{-4}$ |
| 26 | $\mathrm{J}_{\mathrm{C}, 4}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0 | 0 | $4 \cdot 10^{-4}$ | $4 \cdot 10^{-4}$ |
| 27 | $\mathrm{J}_{\mathrm{D}, 4}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0 | 0 | $4 \cdot 10^{-4}$ | $4 \cdot 10^{-4}$ |
| 28 | $\mathrm{J}_{\mathrm{D}, 2}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0 | 0 | $4 \cdot 10^{-4}$ | $4 \cdot 10^{-4}$ |
| 29 | $\mathrm{J}_{\mathrm{E}, 4}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | 0 | 0 | $4 \cdot 10^{-4}$ | $3.2134 \cdot 10^{-4}$ |

The modifications have been weighed through the following regularization that reduces the cost of modifying the nodal inertias, which have been set to zero in the initial nominal model and thus are more uncertain:

$$
\lambda \boldsymbol{\Omega}=2 \cdot 10^{6}\left[\begin{array}{cc}
\mathbf{I}_{[22 \times 22]} & \mathbf{0}_{[22 \times 7]}  \tag{5.25}\\
\mathbf{0}_{[7 \times 22]} & 10^{-3} \cdot \mathbf{I}_{[7 \times 7]}
\end{array}\right]
$$

### 5.5.4. Experimental results

The effectiveness of the updating has been evaluated out by computing the eigenfrequencies and the mode shapes of the updated model and by comparing them with those of the initial nominal one through the percentage frequency $\operatorname{error} \varepsilon_{i}$ :

$$
\begin{equation*}
\varepsilon_{i}=100 \frac{\omega_{i}-\hat{\omega}_{i}}{\hat{\omega}_{i}} \tag{5.26}
\end{equation*}
$$

and the Modal Assurance Criterion MAC:

$$
\begin{equation*}
M A C=\frac{\left(\boldsymbol{\eta}_{i}^{\mathrm{T}} \hat{\boldsymbol{\eta}}_{i}\right)^{2}}{\left\|\boldsymbol{\eta}_{i}\right\|\left\|\hat{\boldsymbol{\eta}}_{i}\right\|} \tag{5.27}
\end{equation*}
$$

The MAC is the mode shape correlation coefficient and should approach the unity.
In Eqs. (5.26) and (5.27) $\omega_{i}$ and $\boldsymbol{\eta}_{i}$ denote the $i$-th eigenfrequency and the eigenvector estimated through the models $(i=1, \ldots, 7)$, while $\hat{\omega}_{i}$ and $\hat{\boldsymbol{\eta}}_{i}$ the measured ones. The comparison proposed in Table 5.4 shows that model updating has improved the model coherence with the experimental measurements. In particular, the greatest improvement is obtained in the estimation of the eigenfrequencies whose values are significantly missed in the nominal model. Indeed, the percentage frequency error for the updated model is smaller than $1 \%$ in four modes, while the average value of such an index is equal to $2.60 \%$ (compared to the $5.75 \%$ in the original model). As for the $M A C$, a less relevant improvement has been obtained since the mode
shape specifications are accurately met also by the original model. Nonetheless, the average $M A C$ increases from 0.896 to 0.909 after the model updating.

The correctness of the mode shapes is confirmed by Figure 5.7, which compares the experimental mode shapes (depicted on the left side of each subplot, as provided by the LMS software) with those computed through the updated model (on the right). The excellent agreement is evident, thus proving the method effectiveness. A further proof of such a good correlation is given by Figure 5.8, which represents the orthogonality (defined through the canonical scalar product I) between different eigenvectors by means of the $M A C$ matrix. Indeed the off-diagonal terms are almost zero.

Table 5.4. Comparison of the experimental and analytical modal properties.

|  | FREQUENCY | FREQUENCY | $\varepsilon_{i}$ | $\varepsilon_{i}$ | MAC | MAC |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| EXPERIMENTAL | NOMINAL | UPDATED | NOMINAL | UPDATED | NOMINAL | UPDATED |
| FREQUENCY | MODEL | MODEL | MODEL | MODEL | MODEL | MODEL |
| [HZ] | $[H z]$ | $[H Z]$ | $[\%]$ | $[\%]$ | $[-]$ | $[-]$ |
|  |  |  |  |  |  |  |
| 13.39 | 13.72 | 13.29 | 2.46 | 0.74 | 0.822 | 0.865 |
| 43.57 | 43.88 | 43.30 | 0.71 | 0.62 | 0.976 | 0.989 |
| 64.65 | 66.96 | 64.78 | 3.58 | 0.20 | 0.960 | 0.968 |
| 112.99 | 127.25 | 124.24 | 12.62 | 9.95 | 0.820 | 0.829 |
| 138.15 | 147.23 | 142.54 | 6.57 | 3.17 | 0.827 | 0.830 |
| 154.73 | 171.95 | 159.37 | 11.13 | 3.00 | 0.928 | 0.944 |
| 198.56 | 204.93 | 197.43 | 3.21 | 0.57 | 0.933 | 0.939 |



## Experimental model

$\qquad$

Figure 5.7. Mode shapes of the first seven modes: experimental (left) and updated (right)


Figure 5.8. MAC matrix between the experimental modes and the updated model modes

### 5.6. Conclusions

A strategy for updating models of flexible link mechanisms has been proposed and validated. The system dynamics is represented through a nonlinear FE model which expresses the coupling between the small elastic deformations and the large motion. The latter can be represented through the rigid-body motion of an equivalent rigid-link system, named the ERLS, which is assumed as a moving reference for defining the small elastic displacements and which is forced to track the real flexible mechanism. The use of the ERLS simplifies the model formulation and allows for the correct representation of the system dynamics, as long as the inertial and elastic parameters are correctly tuned through model updating.

In order to perform model updating through modal analysis, the nonlinear model is linearized about a stable equilibrium configuration where the experimental modal analysis is performed. Then, the measured eigenvectors are transformed from the physical frame to the moving reference defined by the ERLS, by means of the kinematic constraint equations of the rigid-link mechanism and of the definition of the ERLS. Coordinate expansion is also adopted to meet the size of the model and to estimate the elastic rotations.

The calculation of the update of the model parameters is cast as a constrained inverse eigenvalue problem, whose solution can take advantage of several effective and reliable numerical methods available in commercial software. Constraints on the single updating parameters and on 106
combination of more parameters define a feasible set, which ensures the physical meaning of the computed values. The use of preconditioners and of the regularization also improves the numerical reliability of the convex formulation proposed.

The proposed formulation has some strengths that make it suitable for the updating of these complicate models:

- the method can handle an arbitrary number of eigenpairs and updating parameters;
- the experimental eigenvectors can have any arbitrary normalization;
- constraints assure the physical meaning of the updating;
- no iterations are required.

The method has been applied to a very challenging test: the model of a planar, six-bar linkage mechanism with five flexible links. The model is updated to improve the representation of the 7 lowest frequency modes, within the range from 0 to 200 Hz . The experimental results obtained are highly satisfactory, given the complexity of the system investigated, the simultaneous presence of more eigenpairs to be matched, the unavoidable uncertainty in performing accurate measurements in a mechanism (which has "rigid-body" motion) and the presence of constraints bounding the parameter updating. Overall, a meaningful improvement in the eigenvector and eigenfrequency estimation has been obtained.

## Chapter 6. Partial pole placement

### 6.1. Introduction

The eigenvalues of a vibrating system, which are usually called poles in control theory, define the frequency of the dynamic response, the damping ratio and the settling time. In fact, pole placement is a widely employed technique for vibration control. The aim of pole placement is to properly design a controller in such a way that the closed-loop poles are exactly the desired complex numbers. In particular, in this Chapter it is addressed the case in which only one independent control force is exerted on the system, that is known as rank-1 control. In practice, the controlled system equations are

$$
\left\{\begin{array}{l}
\mathbf{M} \ddot{\boldsymbol{q}}+\mathbf{C} \dot{\boldsymbol{q}}+\mathbf{K} \boldsymbol{q}=\boldsymbol{b} u  \tag{6.1}\\
u=-\boldsymbol{f}^{T} \dot{\boldsymbol{q}}-\boldsymbol{g}^{T} \boldsymbol{q}
\end{array}\right.
$$

where $\mathbf{M}, \mathbf{C}, \mathbf{K} \in \mathbb{R}^{N \times N}$ are, respectively, the mass, damping and stiffness matrices, $\boldsymbol{q} \in \mathbb{R}^{N}$ is the displacements vector, $\boldsymbol{b} \in \mathbb{R}^{N}$ represents the distribution of the control forces and $\boldsymbol{f}, \boldsymbol{g} \in \mathbb{R}^{N}$ are the control gains. Given a self-conjugate set of desired poles $\mu_{1}, \ldots, \mu_{2 N} \in \mathbb{C}$, the aim of pole placement is to find such values for $\boldsymbol{f}$ and $\boldsymbol{g}$ that the closed-loop poles are exactly the prescribed ones. It can be proved that a solution exists if and only if

$$
\begin{equation*}
\operatorname{rank}\left[\lambda_{i}^{2} \mathbf{M}+\lambda_{i} \mathbf{C}+\mathbf{K} \mid \boldsymbol{b}\right]=N \tag{6.2}
\end{equation*}
$$

for all the open-loop poles $\lambda_{i}, i=1, \ldots, 2 N$. A system that satisfies this condition is said to be controllable.

For what concerns gain computation, many methods have been proposed in the literature, thus far. A groundbreaking advancement in this field is constituted by the method [36] that relies only on the receptances of the system, rather than on the usual state-space model. The receptance formulation has been later applied to suppression of friction-induced vibration [87], robust pole placement [88], control of flexible link mechanisms [89] and more.

Another significant progress is represented by the methods for partial pole placement, that consist in the assignment of an incomplete set of poles, that is $\mu_{1}, \ldots, \mu_{p}$ with $p<2 N$. Such a necessity is motivated by the fact that often only few dominant poles characterize the dynamics of the system, thus it is advantageous to modify only the most important poles. In such a way, the order of the problem can be considerably reduced, with a great benefit in numerical computation, which is critical for large-scale systems. The difficulty in partial pole placement relies in the fact that if the assignment of $p$ poles changes appreciably the remaining $2 N-p$ ones, the closed-loop system can behave unexpectedly, in some cases even exhibiting instability. Typically, this issue is overcome using conditions that assure that the $2 N-p$ openloop poles are retained in the closed-loop configuration.

In some applications, however, the exact location of the poles is not a critical issue, as long as poles are contained in suitable regions of the complex plane. Indeed, active control has been also employed for such a task, that is usually referred to as regional pole placement. Notable examples of this approach are [90, 91], in which Linear Matrix Inequalities are exploited. Such a mathematical formulation allows to conveniently handle rather complicated specifications on poles location.

In the present Chapter it is proposed a new method for control gain synthesis that enables to perform exact pole placement of $p$ poles, and regional pole placement of the remaining $2 N-p$ ones. This problem has been solved for the control of power systems in [92, 93], but such a method employs mathematical tools which are not numerically reliable (e.g. the characteristic polynomial, the controllability matrix) and hence are not suitable for illconditioned vibrating systems represented through medium or large dimensional matrices with high condition number.

In this work, a new receptance based method for partial pole placement and a state-space method for regional pole placement are combined. The capabilities of the proposed approach are assessed in the design of a controller for a four-bar flexible link mechanism of the type described in Chapter 5.

### 6.2. Description of the method

Let us consider the controlled system (6.1) and suppose that it is wanted to find such control gains $\boldsymbol{f}, \boldsymbol{g}$ that the complex scalars $\mu_{1}, \ldots, \mu_{p}$, with $p \leq 2 N$ are closed-loop poles. The case of complete pole placement, i.e. $p=2 N$, is addressed in [36], where it is defined

$$
\begin{equation*}
\boldsymbol{t}_{i}=\mathbf{H}\left(\mu_{i}\right) \boldsymbol{b} \tag{6.3}
\end{equation*}
$$

with $\mathbf{H}(\lambda)=\left[\lambda^{2} \mathbf{M}+\lambda \mathbf{C}+\mathbf{K}\right]^{-1}$ being the receptance matrix, and then it is proposed to solve the linear system

$$
\left[\begin{array}{cc}
\mu_{1} \boldsymbol{t}_{1}^{T} & \boldsymbol{t}_{1}^{T}  \tag{6.4}\\
\mu_{2} \boldsymbol{t}_{2}^{T} & \boldsymbol{t}_{2}^{T} \\
\vdots & \vdots \\
\mu_{p} \boldsymbol{t}_{p}^{T} & \boldsymbol{t}_{p}^{T}
\end{array}\right]\left\{\begin{array}{c}
\boldsymbol{f} \\
\boldsymbol{g}
\end{array}\right\}=\left\{\begin{array}{c}
-1 \\
-1 \\
\vdots \\
-1
\end{array}\right\} .
$$

If $p=2 N$, the matrix

$$
\mathbf{T}=\left[\begin{array}{cc}
\mu_{1} \boldsymbol{t}_{1}^{T} & \boldsymbol{t}_{1}^{T}  \tag{6.5}\\
\mu_{2} \boldsymbol{t}_{2}^{T} & \boldsymbol{t}_{2}^{T} \\
\vdots & \vdots \\
\mu_{p} \boldsymbol{t}_{p}^{T} & \boldsymbol{t}_{p}^{T}
\end{array}\right]
$$

is a square matrix and, under mild assumptions specified in [36], it is invertible, hence the solution of (6.4) exists, it is real and unique.

If $p<2 N$, the matrix $\mathbf{T}$ is rectangular and the system (6.4) is underdetermined. Without loss of generality, it can be assumed that $\operatorname{rank}(\mathbf{T})=p$ thus the solutions can be represented as

$$
\left\{\begin{array}{l}
f  \tag{6.6}\\
g
\end{array}\right\}=\left\{\begin{array}{l}
f_{0} \\
g_{0}
\end{array}\right\}+\mathbf{V} \boldsymbol{k}_{r}
$$

where $\left\{\boldsymbol{f}_{0}^{T} \boldsymbol{g}_{0}^{T}\right\}^{T} \in \mathbb{R}^{2 N}$ is a particular solution of (6.4), $\mathbf{V}$ is a $2 N \times(2 N-p)$ matrix whose columns span $\operatorname{ker}(\mathbf{T})$ and $\boldsymbol{k}_{r} \in \mathbb{R}^{2 N-p}$ is a vector of parameters that are to be determined. The particular solution $\left\{\boldsymbol{f}_{0}^{T} \boldsymbol{g}_{0}^{T}\right\}^{T}$ can be computed either solving a least-squares problem, or
arbitrarily adding $2 N-p$ rows to $\mathbf{T}$, corresponding to as many poles, which must be judiciously chosen to make the matrix invertible. In the first case, the explicit solution is

$$
\left\{\begin{array}{l}
f_{0}  \tag{6.7}\\
\boldsymbol{g}_{0}
\end{array}\right\}=\left(\mathbf{T}^{T} \mathbf{T}\right)^{-1} \mathbf{T}^{T}\left\{\begin{array}{c}
-1 \\
-1 \\
\vdots \\
-1
\end{array}\right\}
$$

Regardless of how $\left\{\boldsymbol{f}_{0}^{T} \boldsymbol{g}_{0}^{T}\right\}^{T}$ is computed, it follows from [36] that for any choice of $\boldsymbol{k}_{r} \in \mathbb{R}^{2 N-p}$, the control gains $\boldsymbol{f}, \boldsymbol{g}$ computed as in (6.6) lead to a closed-loop system featuring the poles $\mu_{1}, \ldots, \mu_{p}$.

The aim of the following sections is to cast the condition of the remaining $2 N-p$ poles to belong to certain regions of the complex plane in a computationally sound way and to explain how to choose properly the parameter vector $\boldsymbol{k}_{r}$ to assure that such a condition is fulfilled.

### 6.2.1. Linear Matrix Inequalities

The most suitable mathematical tool for regional pole placement is constituted by Linear Matrix Inequalities (or LMIs). In order to exploit LMIs it is necessary to recast the system model in a first-order state-space form, namely

$$
\begin{equation*}
\dot{\boldsymbol{x}}=\mathbf{A x}+\mathbf{B} u \tag{6.8}
\end{equation*}
$$

The chief advantage of the LMIs is that they allow a characterization of the desired constraints that is linear in the state matrix $\mathbf{A}$. Other representations, such as [94], are polynomial and hence not necessarily convex. A brief explanation of Linear Matrix Inequality theory is provided in the present Section, whereas for the details it is suggested to refer to [91].

Let $\mathcal{D}$ be a subset of the complex plane. The system (6.8) is said to be $\mathcal{D}$-stable if all its poles belong to $\mathcal{D}$. There is a class of regions for which there is a simple characterization of $\mathcal{D}$ stability: the LMI regions. $\mathcal{D} \subseteq \mathbb{C}$ is a LMI region if there exist a symmetric matrix $\mathbf{R} \in \mathbb{R}^{m \times m}$ and a matrix $\mathbf{S} \in \mathbb{R}^{m \times m}$ such that

$$
\begin{equation*}
\mathcal{D}=\left\{z \in \mathbb{C}: \quad f_{\mathcal{D}}(z) \prec 0\right\} \tag{6.9}
\end{equation*}
$$

where $f_{\mathcal{D}}(z)=\mathbf{R}+z \mathbf{S}+\bar{z} \mathbf{S}^{T}$, and $\prec 0$ means negative definite. It is easy to prove that LMI regions are convex and symmetric with respect to the real axis. LMIs allow to represent a very large class of subsets of the complex plane, that includes conic sectors, vertical half-planes, horizontal strips, ellipses, parabolas, hyperbolic sectors and intersections of the above. For example, the left half-plane can be proved to be a LMI region simply setting $\mathbf{R}=0$ and $\mathbf{S}=1$. Other examples of LMI regions are shown in Figure 6.1.

The following theorem, whose proof can be found in [91], characterizes $\mathcal{D}$-stability for LMI regions.

Theorem. Let $\mathcal{D}$ be a LMI region. The system $\dot{\boldsymbol{x}}=\mathbf{A x}$ is $\mathcal{D}$-stable if and only if there exists a symmetric matrix $\mathbf{X} \succ 0$ such that

$$
\begin{equation*}
\mathbf{R} \otimes \mathbf{X}+\mathbf{S} \otimes(\mathbf{A X})+\mathbf{S}^{T} \otimes(\mathbf{A X})^{T} \prec 0 \tag{6.10}
\end{equation*}
$$

where $\otimes$ is the Kronecker product and $\mathbf{R}, \mathbf{S}$ are the matrices that appear in the LMI characterization of $\mathcal{D}$.

This result generalizes the Lyapunov stability, which is in fact a special case of Eq. (6.10), for $\alpha=0$ and $\beta=1$. From the computational point of view, the theorem provides a practical way to express a constraint on poles location, since positive (or negative) definite conditions such as (6.10) define convex sets.

maximum natural frequency $\omega<r$

minimum damping ratio $\zeta>\cos \theta$

minimum decay rate $\zeta \omega>\alpha$

Figure 6.1. LMI regions.

### 6.2.2. Control gain synthesis

The actual computation of the control gains $\boldsymbol{f}, \boldsymbol{g}$ that comply with the desired specifications is the results of three stages: first of all, a state-space representation of the controlled system is employed, then the system is reduced to take into account only the $2 N-p$ poles whose location is not exactly imposed, and finally the gains are computed solving an optimization problem with LMI constraints.

The state-space representation that is suggested for the purpose is the very common first-order model

$$
\left\{\begin{array}{l}
\dot{\boldsymbol{x}}=\mathbf{A} \boldsymbol{x}+\mathbf{B} u  \tag{6.11}\\
u=-\boldsymbol{k}^{T} \boldsymbol{x}
\end{array}\right.
$$

where

$$
x=\left\{\begin{array}{l}
\dot{q} \\
q
\end{array}\right\} \quad \mathbf{A}=\left[\begin{array}{cc}
-\mathbf{M}^{-1} \mathbf{C} & -\mathbf{M}^{-1} \mathbf{K} \\
\mathbf{I} & \mathbf{0}
\end{array}\right] \quad \mathbf{B}=\left\{\begin{array}{c}
\mathbf{M}^{-1} \boldsymbol{b} \\
\mathbf{0}
\end{array}\right\} \quad \boldsymbol{k}=\left\{\begin{array}{l}
\boldsymbol{f} \\
\boldsymbol{g}
\end{array}\right\} .
$$

In order to achieve the desired reduction, it is employed the non-invertible transformation $\boldsymbol{x}=\mathbf{V} \hat{\boldsymbol{x}}$, where $\hat{\boldsymbol{x}} \in \mathbb{R}^{2 N-p}$ is the state of the reduced system and $\mathbf{V}$ is defined in (6.6). Let $\mathbf{U} \in \mathbb{R}^{(2 N-p) \times 2 N}$ be the Moore-Penrose pseudoinverse of $\mathbf{V}$, then since $\mathbf{V}$ has linearly independent columns, $\mathbf{U}$ is its left inverse, that is $\mathbf{U V}=\mathbf{I}_{2 N-p}$. Hence the reduced system is

$$
\left\{\begin{array}{l}
\dot{\hat{\boldsymbol{x}}}=\mathbf{U A V} \hat{\boldsymbol{x}}+\mathbf{U B} \hat{u}  \tag{6.12}\\
\hat{u}=-\boldsymbol{k}^{T} \mathbf{V} \hat{\boldsymbol{x}} .
\end{array}\right.
$$

The dynamics of the closed-loop system is described by the equation

$$
\begin{equation*}
\dot{\hat{\boldsymbol{x}}}=\mathbf{U A V} \hat{\boldsymbol{x}}-\mathbf{U B} \boldsymbol{k}^{T} \mathbf{V} \hat{\boldsymbol{x}} \tag{6.13}
\end{equation*}
$$

and recalling that the control gains are wanted to satisfy (6.6), that is $\boldsymbol{k}=\boldsymbol{k}_{0}+\mathbf{V} \boldsymbol{k}_{r}$, where $\boldsymbol{k}_{0}=\left\{\boldsymbol{f}_{0}^{T} \boldsymbol{g}_{0}^{T}\right\}^{T}$, the new control system becomes

$$
\left\{\begin{array}{l}
\dot{\hat{\boldsymbol{x}}}=\mathbf{U}\left(\mathbf{A}-\mathbf{B} \boldsymbol{k}_{0}\right) \mathbf{V} \hat{\boldsymbol{x}}+\mathbf{U B} \hat{u}  \tag{6.14}\\
\hat{u}=-\left(\mathbf{V}^{T} \mathbf{V} \boldsymbol{k}_{r}\right)^{T} \hat{\boldsymbol{x}} .
\end{array}\right.
$$

For clarity, it is defined $\hat{\mathbf{A}}=\mathbf{U}\left(\mathbf{A}-\mathbf{B} \boldsymbol{k}_{0}\right) \mathbf{V}, \hat{\mathbf{B}}=\mathbf{U B}$ and $\hat{\boldsymbol{k}}=\mathbf{V}^{T} \mathbf{V} \boldsymbol{k}_{r}$ in such a way that the system can be written as

$$
\left\{\begin{array}{l}
\dot{\hat{\boldsymbol{x}}}=\hat{\mathbf{A}} \hat{\boldsymbol{x}}+\hat{\mathbf{B}} \hat{u}  \tag{6.15}\\
\hat{u}=-\hat{\boldsymbol{k}}^{T} \hat{\boldsymbol{x}}
\end{array}\right.
$$

The remarkable property of this system is that its closed-loop poles are exactly the $2 N-p$ non-assigned closed-loop poles of (6.11). This property will be further analysed in the next Section.

The goal is to find a control gain $\hat{\boldsymbol{k}}$ such that the closed-loop poles of (6.15) lie in a prescribed LMI region $\mathcal{D}$. In practice, it is wanted that Eq. (6.10) holds for $\hat{\mathbf{A}}_{\mathrm{CL}}=\hat{\mathbf{A}}-\hat{\mathbf{B}} \hat{\boldsymbol{k}}^{T}$. Substituting, the matrix inequality becomes non-linear, in fact both $\hat{\boldsymbol{k}}$ and $\mathbf{X}$ are unknown and hence the following

$$
\begin{equation*}
\mathbf{R} \otimes \mathbf{X}+\mathbf{S} \otimes\left(\hat{\mathbf{A}} \mathbf{X}-\hat{\mathbf{B}} \hat{\boldsymbol{k}}^{T} \mathbf{X}\right)+\mathbf{S}^{T} \otimes\left(\hat{\mathbf{A}} \mathbf{X}-\hat{\mathbf{B}} \hat{\boldsymbol{k}}^{T} \mathbf{X}\right)^{T} \tag{6.16}
\end{equation*}
$$

contains bilinear unknown terms. Linearity is easily restored introducing the auxiliary variables $\boldsymbol{l}=\hat{\boldsymbol{k}}^{T} \mathbf{X}$, thus leading to the Linear Matrix Inequality

$$
\begin{equation*}
\mathcal{F}(\mathbf{X}, \boldsymbol{l})=\mathbf{R} \otimes \mathbf{X}+\mathbf{S} \otimes(\hat{\mathbf{A}} \mathbf{X}-\hat{\mathbf{B}} \boldsymbol{l})+\mathbf{S}^{T} \otimes(\hat{\mathbf{A}} \mathbf{X}-\hat{\mathbf{B}} \boldsymbol{l})^{T} \prec 0 \tag{6.17}
\end{equation*}
$$

It is suggested to compute the control gains solving the following convex optimization problem

$$
\begin{equation*}
\min _{\substack{\mathcal{F}(\mathbf{X}, l)<0 \\ \mathbf{x} \succ 0}}\|l\|^{2} \tag{6.18}
\end{equation*}
$$

and then computing the desired $\hat{\boldsymbol{k}}$ from the relation $\hat{\boldsymbol{k}}^{T}=\boldsymbol{l} \mathbf{X}^{-1}$. The knowledge of $\hat{\boldsymbol{k}}$ enables to obtain $\boldsymbol{k}_{r}$ and, finally, $\boldsymbol{f}$ and $\boldsymbol{g}$ from (6.6). A controlled system with such gains is assured to feature $\mu_{1}, \ldots, \mu_{p}$ as closed-loop poles, and the other ones are proved to belong to $\mathcal{D}$. Moreover, it is worth to notice that

$$
\begin{equation*}
\|\boldsymbol{l}\|=\left\|\hat{\boldsymbol{k}}^{T} \mathbf{X}\right\|=\|\mathbf{X} \hat{\boldsymbol{k}}\| \tag{6.19}
\end{equation*}
$$

and, since $\mathbf{X}$ is a positive definite matrix, the minimization of the norm of $\boldsymbol{l}$ allows to keep the controller effort small.

### 6.2.3. Details on the reduced system model

The aim of the present Section is to give a deeper insight on the fact that the non-assigned poles of (6.11) are indeed the poles of the reduced system (6.15).

In [95] it is introduced the notion of inclusion in dynamic systems. Although such a notion does not apply to the systems of interest, the proof in the cited paper gives a hint on a possible way to prove that the poles of the full-order system are exactly the $p$ desired poles $\mu_{1}, \ldots, \mu_{p}$ plus the $2 n-p$ poles of the reduced system.

Let us denote $\tilde{\mathbf{V}}$ a matrix whose columns span $\operatorname{ker}(\mathbf{U})$, and $\tilde{\mathbf{U}}$ the unique left inverse of $\tilde{\mathbf{V}}$ such that $\operatorname{ker}(\tilde{\mathbf{U}})=\operatorname{im}(\tilde{\mathbf{V}})$. It is also defined $\mathbf{W}=[\mathbf{V} \tilde{\mathbf{V}}]$ and $\mathbf{Z}=\mathbf{A}_{\mathrm{CL}}-\mathbf{V} \hat{\mathbf{A}}_{\mathrm{CL}} \mathbf{U}$. It can be proved that

$$
\mathbf{W}^{-1} \mathbf{A}_{\mathrm{CL}} \mathbf{W}=\left[\begin{array}{cc}
\hat{\mathbf{A}}_{\mathrm{CL}} & \mathbf{U} \mathbf{Z} \tilde{\mathbf{V}}  \tag{6.20}\\
\tilde{\mathbf{U}} \mathbf{Z} \mathbf{V} & \tilde{\mathbf{U}} \mathbf{Z} \tilde{\mathbf{V}}
\end{array}\right],
$$

hence the closed-loop full order system has the same poles as the right hand side of (6.20). For the latter matrix, it is particularly easy to compute the poles, in fact it has been numerically verified that $\mathbf{U Z} \tilde{\mathbf{V}}=\mathbf{0}$. The general proof of this fact is still missing, but it has been verified in all the numerical tests performed, as well as the fact that the eigenvalues of $\tilde{\mathbf{U}} \mathbf{Z} \tilde{\mathbf{V}}$ are exactly $\mu_{1}, \ldots, \mu_{p}$.

### 6.2.4. Extension to asymmetric systems

Even though the linear system (6.4) has been introduced in [36] for symmetric systems only, such a limitation has been overcome in [89] and thus the method can be also employed for the control of systems that feature non-symmetric damping and stiffness matrices, such as for example flexible link mechanisms.

The linear system (6.4) can be adjusted as follows. Let us decompose the damping and stiffness matrices into the sum of the symmetric parts $\mathbf{C}_{s}, \mathbf{K}_{s} \in \mathbb{R}^{N \times N}$ and asymmetric ones $\mathbf{C}_{a}, \mathbf{K}_{a} \in \mathbb{R}^{N \times N}$, namely:

$$
\begin{equation*}
\mathbf{C}=\mathbf{C}_{s}+\mathbf{C}_{a}, \quad \mathbf{K}=\mathbf{K}_{s}+\mathbf{K}_{a} . \tag{6.21}
\end{equation*}
$$

Such a decomposition is not unique and it must be chosen carefully. Thus it is defined $\mathbf{H}_{S}(\lambda)=\left[\lambda^{2} \mathbf{M}+\lambda \mathbf{C}_{s}+\mathbf{K}_{s}\right]^{-1}$ to be the receptance matrix of the open-loop symmetric system, with $h_{N j}^{s}(\lambda)$ being its entry at the $N$-th row and $j$-th column. $k_{a_{j}}$ is the entry at the $j$-th row and $N$-th column of $\mathbf{K}_{a}$ and $c_{a_{j}}$ is the entry at the $j$-th row and $N$-th column of $\mathbf{C}_{a}$. Moreover, $\boldsymbol{e}_{N} \in \mathbb{R}^{N}$ denotes the vector whose $N$-th entry is 1 , while the others are 0 , and $\overline{\boldsymbol{e}}_{N}(\lambda) \in \mathbb{R}^{N}$ is the vector whose $j$-th entry is $k_{a_{j}}+\lambda c_{a_{j}}$, for $j=1, \ldots, N-1$ and the $N$-th entry is 0 . Finally, the definition (6.3) is adjusted as follows

$$
\begin{equation*}
\boldsymbol{t}_{i}=\left(1+\sum_{j=1}^{N-1}\left(k_{a_{j}}+\mu_{i} c_{a_{j}}\right) h_{N j}^{s}\left(\mu_{i}\right)\right) \mathbf{H}_{S}\left(\mu_{i}\right) \boldsymbol{b}-\mathbf{H}_{S}\left(\mu_{i}\right) \overline{\boldsymbol{e}}_{N} \boldsymbol{b}^{T} \mathbf{H}_{S}\left(\mu_{i}\right) \boldsymbol{e}_{N} \tag{6.22}
\end{equation*}
$$

and the linear system (6.4) becomes

$$
\left[\begin{array}{cc}
\mu_{1} \boldsymbol{t}_{1}^{T} & \boldsymbol{t}_{1}^{T}  \tag{6.23}\\
\mu_{2} \boldsymbol{t}_{2}^{T} & \boldsymbol{t}_{2}^{T} \\
\vdots & \vdots \\
\mu_{p} \boldsymbol{t}_{p}^{T} & \boldsymbol{t}_{p}^{T}
\end{array}\right]\left\{\begin{array}{l}
\boldsymbol{f} \\
\boldsymbol{g}
\end{array}\right\}=-\left\{\begin{array}{c}
1+\sum_{j=1}^{N-1}\left(k_{a_{j}}+\mu_{1} c_{a_{j}}\right) h_{N j}^{s}\left(\mu_{1}\right) \\
1+\sum_{j=1}^{N-1}\left(k_{a_{j}}+\mu_{2} c_{a_{j}}\right) h_{N j}^{s}\left(\mu_{2}\right) \\
\vdots \\
1+\sum_{j=1}^{N-1}\left(k_{a_{j}}+\mu_{p} c_{a_{j}}\right) h_{N j}^{s}\left(\mu_{p}\right)
\end{array}\right\} .
$$

The detailed calculations that lead to these expressions can be found in [89]. Nevertheless, the tests carried out, such as the one shown in the next Section, prove that if the definition of $\mathbf{T}$ in Eq. (6.5) is suitably updated, the proposed method allows to effectively perform partial pole placement for asymmetric systems.

### 6.3. Numerical assessment

The method has been validated by synthesizing a controller for the four-bar flexible-link mechanism in Figure 6.2. The links are assumed to be straight and slender steel bars and are modelled through Euler-Bernoulli beam elements. The system model is analogous to the one of the mechanism in Chapter 5. The geometric and inertial characteristics of the mechanism are summarized in Table 6.1. The finite element representation employed for the mechanism, which is illustrated in Figure 6.3, leads to a 16-dof dynamic model ( $N=16$ ), with one rigid-body degree


Figure 6.2. Picture of the mechanism.


Figure 6.3. Finite element model of the four-bar linkage.

Table 6.1. Parameters of the four-bar linkage.

| Parameter | Value |
| :--- | :--- |
| Young's modulus | $210 \cdot 10^{9}[\mathrm{~Pa}]$ |
| Flexural inertia moment | $11.102 \cdot 10^{-10}\left[\mathrm{~m}^{4}\right]$ |
| Beams width | $6 \cdot 10^{-3}[\mathrm{~m}]$ |
| Beams thickness | $6 \cdot 10^{-3}[\mathrm{~m}]$ |
| Mass/unit of length of links | $272 \cdot 10^{-3}[\mathrm{~kg} / \mathrm{m}]$ |
| Crank (link 1) length | $0.360[\mathrm{~m}]$ |
| Coupler (link 2) length | $0.390[\mathrm{~m}]$ |
| Follower (link 3) length | $0.535[\mathrm{~m}]$ |
| Frame (link 4) length | $0.632[\mathrm{~m}]$ |
| Rayleigh damping constants | $\alpha=8.5 \cdot 10^{-2}\left[\mathrm{~s}^{-1}\right]$ |
|  | $\beta=2 \cdot 10^{-5}[\mathrm{~s}]$ |
| Joint masses: | 0 |
| - Joint A | $70 \cdot 10^{-3}[\mathrm{~kg}]$ |
| - Joint B | $70 \cdot 10^{-3}[\mathrm{~kg}]$ |
| - Joint C | 0 |
| - Joint D |  |
| Joint inertia: | $5 \cdot 10^{-4}\left[\mathrm{kgm}^{2}\right]$ |
| - Joint A | 0 |
| - Joint B | 0 |
| - Joint C | $12 \cdot 10^{-6}\left[\mathrm{kgm}{ }^{2}\right]$ |
| - Joint D |  |

of freedom and 15 vibrational modes. A single motor actuates the mechanism, driving the rotation of the crank. This ensures the controllability of the whole 16 -dof system.

Finally, the mechanism is supposed to move on a vertical plane, and so the open-loop system is unstable as a consequence of the presence of gravity. Further details on the system model can be found in [79, 89].

The pole placement problem has been solved in a neighbourhood of the crank horizontal configuration, and it consists in the assignment of 22 poles, while the remaining 10 poles are constrained into a LMI region.

The 22 poles that are exactly assigned are shown in the second column of Table 6.2. The remaining 10 poles are required to belong to the intersection of a conic sector $\mathcal{D}_{1}$ of magnitude $\frac{5}{6} \pi$ (which means that the damping ratio of the associated modes is at least $\zeta_{\text {min }}=0.25$ ) and the left half-plane $\mathcal{D}_{2}$ defined by $\operatorname{Re}(z)<-1000$ (which forces the settling time of such modes to be less than $T_{\text {max }}=0.003 \mathrm{~s}$ ). It is easy to prove that such sets are LMI regions, in fact they can both be represented as in (6.9), with

$$
\begin{aligned}
& f_{\mathcal{D}_{1}}(z)=\left[\begin{array}{cc}
\sin \frac{\vartheta}{2} & -\cos \frac{\vartheta}{2} \\
\cos \frac{v}{2} & \sin \frac{v}{2}
\end{array}\right] z+\left[\begin{array}{cc}
\sin \frac{\vartheta}{2} & \cos \frac{\vartheta}{2} \\
-\cos \frac{\vartheta}{2} & \sin \frac{v}{2}
\end{array}\right] \bar{z}, \quad \vartheta=\frac{5}{6} \pi, \\
& f_{\mathcal{D}_{2}}(z)=2 \cdot 1000+z+\bar{z} .
\end{aligned}
$$

The resulting LMI constraints are defined as in (6.17) and are denoted $\mathcal{F}_{1}$ and $\mathcal{F}_{2}$. Rather than imposing $\mathcal{F}_{1}\left(\mathbf{X}_{1}, \boldsymbol{l}_{1}\right) \prec 0$ for some $\mathbf{X}_{1} \succ 0$ and $\mathcal{F}_{2}\left(\mathbf{X}_{2}, \boldsymbol{l}_{2}\right) \prec 0$ for some $\mathbf{X}_{2} \succ 0$ separately, it is convenient to seek a common solution $\mathbf{X}=\mathbf{X}_{1}=\mathbf{X}_{2}$ and $\boldsymbol{l}=\boldsymbol{l}_{1}=\boldsymbol{l}_{2}$, which is more conservative but allows to reduce the number of involved unknowns.

The poles obtained with the proposed partial pole placement method are shown in the third column of Table 6.2 , which proves that the 22 assigned poles are accurately matched. For clarity, not all the closed-loop poles are represented in Figure 6.4, but the plot allows to verify that the remaining poles actually belong to the prescribed LMI region $\mathcal{D}_{1} \cap \mathcal{D}_{2}$.

Table 6.2. Pole comparison.

| Open-loop poles | Desired poles | Closed-loop poles |
| :---: | :---: | :---: |
| +3.16 | -15 | -15.00+0.00i |
| -3.16 | -20 | $-20.00+0.00 \mathrm{i}$ |
| $-1.2+226.2 \mathrm{i}$ | -10-220i | -9.36-220.78i |
| -1.2-226.2i | $-10+220 \mathrm{i}$ | $-9.36+220.78 \mathrm{i}$ |
| $-2.4+322.2 \mathrm{i}$ | -15-320i | -14.74-321.31i |
| -2.4-322.2i | $-15+320 \mathrm{i}$ | $-14.74+321.31 \mathrm{i}$ |
| -5.3+483.1i | -50-500i | -63.82-505.74i |
| -5.3-483.1i | -50+500i | $-63.82+505.74 \mathrm{i}$ |
| $-20.9+958.4 \mathrm{i}$ | -100-900i | -100.24-897.02i |
| -20.9-958.4i | $-100+900 i$ | $-100.24+897.02 \mathrm{i}$ |
| -46.2+1427.4i | -50-1400i | -49.54-1399.46i |
| -46.2-1427.4i | -50+1400i | -49.54+1399.46i |
| $-58.5+1605.2 \mathrm{i}$ | -100-1600i | -96.32-1597.37i |
| -58.5-1605.2i | $-100+1600 \mathrm{i}$ | $-96.32+1597.37 \mathrm{i}$ |
| $-121.1+2308.9 \mathrm{i}$ | -130-2300i | -129.95-2300.00i |
| -121.1-2308.9i | $-130+2300 i$ | $-129.95+2300.00 \mathrm{i}$ |
| -287.6+3551.7i | -290-3550i | -289.99-3549.98i |
| -287.6-3551.7i | $-290+3550 \mathrm{i}$ | $-289.99+3549.98 \mathrm{i}$ |
| $-390.7+4134.8 i$ | -400-4000i | -400.00-4000.08i |
| -390.7-4134.8i | $-400+4000 \mathrm{i}$ | -400.00+4000.08i |
| -946.9+6396.1i | -1000-6300i | -999.99-6300.05i |
| -946.9-6396.1i | $-1000+6300 i$ | -999.99+6300.05i |
| $-1819.9+8777.4 \mathrm{i}$ |  | $-1330.83+0.00 \mathrm{i}$ |
| -1819.9-8777.4i |  | -2202.36-7980.76i |
| $-3440.7+11834.8 i$ |  | $-2202.36+7980.76 i$ |
| -3440.7-11834.8i |  | -4488.89-13453.89i |
| $-6348.2+15490.7 \mathrm{i}$ |  | -4488.89+13453.89i |
| -6348.2-15490.7i |  | -22838.21-22166.55i |
| -23417.4+22032.6i |  | $-22838.21+22166.55 i$ |
| -23417.4-22032.6i |  | -42526.20-5837.17i |
| -40747.3+11773.9i |  | $-42526.20+5837.17 \mathrm{i}$ |
| -40747.3-11773.9i |  | -418716.84+0.00i |



Figure 6.4. Pole map: low frequency detail.

### 6.4. Conclusion

Partial pole placement is a common task in the design of vibrating systems, because it is often necessary to set the location only of the most relevant poles. The remaining poles can often constitute an issue, if not dealt with properly, because spill-over of such poles can lead to instability. Techniques that allow to prevent spill-over are available in the literature, but in this work it is rather proposed to constrain the non-assigned poles into desired regions of the complex plane. The greater freedom in the location of the poles can be exploited to minimize the norm of the control gains, and other properties to be optimized are under investigation. The method has been successfully employed to synthesize the gains for the controller of a four-bar flexible link mechanisms, thus showing the capabilities of the proposed approach.

## Chapter 7. Hybrid control

### 7.1. Introduction

Let us consider the $N$ degrees of freedom undamped vibrating system which can be controlled exerting $N_{B}$ control forces, namely:

$$
\begin{equation*}
\mathbf{M} \ddot{\boldsymbol{q}}(t)+\mathbf{K} \boldsymbol{q}(t)=\mathbf{B} \boldsymbol{v}(t) . \tag{7.1}
\end{equation*}
$$

where $\mathbf{M}, \mathbf{K} \in \mathbb{R}^{N \times N}$ are respectively the symmetric mass and stiffness matrices, $\boldsymbol{q} \in \mathbb{R}^{N}$ is the displacement vector, $\mathbf{B} \in \mathbb{R}^{N \times N_{B}}$ represents the distribution of $N_{B}$ independent control forces and $\boldsymbol{v}(t) \in \mathbb{R}^{N_{B}}$ is the control force vector.

Although such an approach requires the use of actuators and sensors, it is often employed in eigenstructure assignment. For what concerns the assignment of just the eigenvalues, i.e. the so-called "pole placement", several methods are available [29, 30, 87]. The assignment of eigenvectors, in contrast, is a much more challenging problem which is not completely solved. In fact, it can be proven that active control can assign only eigenvectors that belong to certain vector spaces which depend on the physical properties of the system and on the characteristics of its actuation [33, 63]. In particular, the dimension of the space of allowable eigenvectors matches the number of independent actuation forces. Thus if the system is highly underactuated it is very unlikely that exact assignment of eigenvectors can be accomplished.

Recently, passive and active control have been combined also to develop hybrid methods, and few promising results have appeared in the literature, so far. For example, Corr and Clark in [96] achieve vibration confinement using mechanical dampers to reduce active power requirements for active control. In a similar way, Tang and Wang in [97] obtain the same objective employing piezoelectric materials both as active actuators and as passive dampers. Diaz et al. in [98], instead, manage to improve the performances of input shaping by modifying the link in a flexible manipulator. More recently, Ouyang in [99] consecutively employed structural modification and feedback control for eigenvalue assignment in asymmetric systems. The general problem of eigenstructure assignment has been instead tackled in [100] by Richiedei and Trevisani, who developed a method for the simultaneous computation of the
structural modifications and control gains synthesis to overcome the inherent limitations active control, for what concerns the assignment of eigenvectors.

Starting from such a work, a different mathematical frame for the problem solution is proposed in this Chapter to handle highly underactuated systems and large dimensional models, by taking advantage of some recent results of matrix analysis. Similarly to [100], the idea consists in exploiting Dynamic Structural Modification to widen the range of the eigenvectors that can be achieved through active control, even when the number of available actuation forces is low. The assignment problem is handled by solving a matrix rank minimization problem. To that end, two different algorithms are discussed, one that relies on non-convex optimization and the other that, instead, is based on semi-definite optimization. The effectiveness of the developed hybrid approach is assessed in two different test-cases. The first one allows to compare the proposed method with the state-of-the art methods in a benchmark problem. The second one, instead, enables to show the peculiarities of the two considered algorithms, with regard to the choice of the actuation.

### 7.2. Problem formulation

### 7.2.1. Eigenstructure assignment through state feedback active control

The aim of this Section is to discuss the limitation of Eigenstructure Assignment through active control. It is supposed that the system model (7.1) is employed, with the following control force:

$$
\begin{equation*}
\boldsymbol{f}_{C}(t)=\mathbf{B} \boldsymbol{v}(t), \tag{7.2}
\end{equation*}
$$

where $\mathbf{B} \in \mathbb{R}^{N \times N_{B}}$ represents the distribution of $N_{B}$ independent control forces and $\boldsymbol{v}(t) \in \mathbb{R}^{N_{B}}$ is the control force vector.

It is also supposed that a position-feedback control law is employed, such as

$$
\begin{equation*}
\boldsymbol{v}(t)=\mathbf{G}^{T}\left(\boldsymbol{q}(t)-\boldsymbol{q}_{\mathrm{ref}}(t)\right) \tag{7.3}
\end{equation*}
$$

where $\mathbf{G} \in \mathbb{R}^{N \times N_{B}}$ is the control gain matrix and $\boldsymbol{q}_{\mathrm{ref}}(t)$ is the position reference.

The closed-loop system is therefore:

$$
\begin{equation*}
\mathbf{M} \ddot{\boldsymbol{q}}(t)+\left(\mathbf{K}-\mathbf{B G}^{T}\right) \boldsymbol{q}(t)=-\mathbf{B G}^{T} \boldsymbol{q}_{\mathrm{ref}}(t) . \tag{7.4}
\end{equation*}
$$

The closed-loop system eigenvalues and eigenvectors are, respectively, the scalars $\omega^{2} \in \mathbb{R}$ and the vectors $\boldsymbol{u} \in \mathbb{R}^{N}$ that solve the following equation:

$$
\begin{equation*}
\left[\mathbf{K}-\mathbf{B G}^{T}-\omega^{2} \mathbf{M}\right] \boldsymbol{u}=\mathbf{0} . \tag{7.5}
\end{equation*}
$$

The eigenfrequency $\omega$ and its associated eigenvector $\boldsymbol{u}$, constitute together an eigenpair. Given a set of $N_{e}$ desired eigenpairs $\left(\omega_{1}, \boldsymbol{u}_{1}\right),\left(\omega_{2}, \boldsymbol{u}_{2}\right), \ldots,\left(\omega_{N_{e}}, \boldsymbol{u}_{N_{e}}\right)$, eigenstructure assignment consists in finding the controller gain matrix $\mathbf{G} \in \mathbb{R}^{N \times N_{B}}$ such that the closed-loop system features the eigenpairs $\left(\omega_{i}, \boldsymbol{u}_{i}\right)$ for $i=1, \ldots, N_{e}$.

### 7.2.2. Existence of solutions for active eigenstructure assignment

The existence of feedback gains that result in the assignment of a prescribed set of eigenvalues is a problem exhaustively investigated in literature and is governed by the notion of system controllability. In particular, a system is said to be controllable if and only if the following condition holds,

$$
\begin{equation*}
\operatorname{rank}\left(\left[\mathbf{K}-\omega_{i}^{2} \mathbf{M} \mid \mathbf{B}\right]\right)=N \tag{7.6}
\end{equation*}
$$

for all the open-loop eigenfrequencies $\omega_{i}(i=1, \ldots, N)$. If the system is controllable then it is possible to place the closed-loop eigenvalues in the desired locations of the complex plane by means of state-feedback control.

The fulfillment of condition (7.6), however, is not sufficient to assure that also the closed-loop eigenvectors can be properly assigned. In fact, it can be proved that for any eigenfrequency $\omega$, the corresponding eigenvector must belong to a particular vector space $\Psi(\omega) \subseteq \mathbb{R}^{N}$. Indeed equation (7.5), that must be satisfied in order to assign the eigenpair ( $\omega, \boldsymbol{u}$ ), can be rearranged as follows:

$$
\begin{equation*}
\left[\mathbf{K}-\omega^{2} \mathbf{M}\right] \boldsymbol{u}=\mathbf{B G}^{T} \boldsymbol{u} \tag{7.7}
\end{equation*}
$$

If equation (7.7) is left multiplied by the open-loop system receptance matrix at $\omega$, i.e. $\mathbf{H}(\omega)=\left[\mathbf{K}-\omega^{2} \mathbf{M}\right]^{-1}$, then it follows that

$$
\begin{equation*}
\boldsymbol{u}=\mathbf{H}(\omega) \mathbf{B} \boldsymbol{x} \tag{7.8}
\end{equation*}
$$

where $\boldsymbol{x}=\mathbf{G}^{T} \boldsymbol{u} \in \mathbb{R}^{N_{B}}$. Given the arbitrariness of the gain matrix $\mathbf{G}$, equation (7.8) implies that any assignable eigenvector $\boldsymbol{u}$ must belong to the vector space spanned by the columns of matrix $[\mathbf{H}(\omega) \mathbf{B}] \in \mathbb{R}^{N \times N_{B}}$. Hence such a space, henceforth denoted the allowable subspace $\Psi(\omega)$, collects the eigenvectors that can be obtained by active control, associated with the eigenfrequency $\omega$. Thus,

$$
\begin{equation*}
\Psi(\omega)=\operatorname{span}\{\mathbf{H}(\omega) \mathbf{B}\} . \tag{7.9}
\end{equation*}
$$

Since both $\mathbf{H}(\omega)$ and $\mathbf{B}$ have maximum rank, the space $\Psi$ is $N_{B}$-dimensional. Therefore, if $N_{B}<N$, namely the system is underactuated, it is not guaranteed that the eigenvector can be exactly assigned. Such a limitation is particularly serious if few control forces are used, that is $N_{B} \ll N$, because the desired eigenvector $\boldsymbol{u}$ is unlikely to lie in the allowable subspace. A workaround that is often used in the literature is to assign the projection of $\boldsymbol{u}$ onto $\Psi$, rather than $\boldsymbol{u}$ itself [33]. In the most critical cases, however, this strategy leads to inadequate closedloop system performances because the projection just roughly approximate $\boldsymbol{u}$. A proposed solution to solve this issue is discussed in the next Section.

### 7.3. Proposed solution

### 7.3.1. Matrix rank minimization problem

By following and improving the idea proposed in [100], it is exploited Dynamic Structural Modification to obtain a new allowable subspace $\hat{\Psi}\left(\omega_{i}\right)$ in such a way that $\boldsymbol{u}_{i} \in \hat{\Psi}\left(\omega_{i}\right)$ for all the $N_{e}$ desired eigenpairs $\left(\omega_{i}, \boldsymbol{u}_{i}\right)$. Compared to such a work, which proposes the first idea of hybrid control to eigenstructure assignment, this one formulates a more effective solution that
takes advantage of some recent theorems in the field of matrix analysis and of the state-of-theart methods for gain computation.

According to the definition of allowable subspace in equation (7.9), $\Psi$ depends on the receptance matrix computed at the desired eigenfrequency, i.e. $\mathbf{H}\left(\omega_{i}\right)$, and on the input matrix B . It will be henceforth assumed that the matrix B cannot be modified, which is a plausible assumption since there are often technical constraints on the number and the location of the actuators. Therefore, an effective way to shape the allowable subspaces is to exploit DSM to modify the receptance matrix for each eigenvalue of interest. In practice, it is wanted to find modification matrices $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$ such that the following relation holds for each assigned eigenpair $\left(\omega_{i}, \boldsymbol{u}_{i}\right)$ :

$$
\begin{equation*}
\boldsymbol{u}_{i} \in \operatorname{span}\left\{\left[(\mathbf{K}+\Delta \mathbf{K})-\omega_{i}^{2}(\mathbf{M}+\Delta \mathbf{M})\right]^{-1} \mathbf{B}\right\}=\hat{\Psi}\left(\omega_{i}\right) . \tag{7.10}
\end{equation*}
$$

The computation of the appropriate modification matrices can be recast as a rank minimization problem. Such a formulation is advantageous because it is not necessary to cope with the receptance matrix and thus no inversions of the matrix collecting the unknowns are needed. Moreover, since rank minimization problems have many applications in different fields [101, $102,103]$, several numerical algorithms have been recently proposed [104, 105, 106].

The proposed formulation is obtained by taking advantage of the well-known Rouché-Capelli theorem of linear algebra. In fact, Eq. (7.10) holds if there are coefficient vectors $\boldsymbol{x}_{i} \in \mathbb{R}^{N_{B}}$ such that

$$
\begin{equation*}
\mathbf{B} \boldsymbol{x}_{i}=\left[(\mathbf{K}+\Delta \mathbf{K})-\omega_{i}^{2}(\mathbf{M}+\Delta \mathbf{M})\right] \boldsymbol{u}_{i} . \tag{7.11}
\end{equation*}
$$

By defining vectors $\boldsymbol{c}_{i}$ as follows,

$$
\begin{equation*}
\boldsymbol{c}_{i}=\left[(\mathbf{K}+\Delta \mathbf{K})-\omega_{i}^{2}(\mathbf{M}+\Delta \mathbf{M})\right] \boldsymbol{u}_{i} \in \mathbb{R}^{N}, \tag{7.12}
\end{equation*}
$$

equations (7.11) can be conveniently written as

$$
\begin{equation*}
\mathbf{B} x_{i}=c_{i} . \tag{7.13}
\end{equation*}
$$

According to the Rouché-Capelli theorem, each one of the systems in Eq. (7.13) can be solved if and only if the rank of the matrix $\mathbf{B}$ (which is assumed to be $N_{B}$ ) is equal to the rank of the augmented matrix $\left[\mathbf{B} \mid \boldsymbol{c}_{i}\right]$. Therefore, the systems in Eq. (7.13) can be solved for all the indices $i=1, \ldots, N_{e}$ if and only if

$$
\begin{equation*}
\operatorname{rank}\left(\left[\mathbf{B}\left|\boldsymbol{c}_{1}\right| \ldots \mid \boldsymbol{c}_{N_{e}}\right]\right)=N_{B} \tag{7.14}
\end{equation*}
$$

In practice, the computation of the suitable modification of the allowable subspaces consists in finding the matrices $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$ such that Eq. (7.14) holds. It is reasonable to restrict the research of the structural modifications to a constraint set $\Gamma \subseteq \mathbb{R}^{N \times N} \times \mathbb{R}^{N \times N}$ to take into account the physical feasibility of the solution, due to the technical or economic constraints on the allowable modifications. In such a way, however, the existence of a solution is not guaranteed. Therefore, it is convenient to use an optimization-based formulation: it is wanted to find the modification matrices $\Delta \mathbf{M}, \Delta \mathbf{K}$ that solve the following problem:

$$
\begin{array}{ll}
\operatorname{minimize} & \operatorname{rank}\left[\mathbf{B}\left|\boldsymbol{c}_{1}\right| \ldots \mid \boldsymbol{c}_{N_{e}}\right]  \tag{7.15}\\
\text { subject to } & (\Delta \mathbf{M}, \Delta \mathbf{K}) \in \Gamma
\end{array}
$$

Since the rank of a matrix is numerically difficult to handle in an optimization problem, the algorithms to solve problem (7.15) must be carefully discussed.

### 7.3.2. Numerical algorithms

Two different numerical algorithms for matrix rank minimization are taken into consideration, by taking advantage of some theoretical results of the linear algebra. The peculiarities of each algorithm will be discussed in the next Sections.

## Algorithm 1

The first algorithm relies on the following theorem of linear algebra [105].

Theorem 1. Let $\mathbf{X} \in \mathbb{R}^{m \times n}$, with arbitrary $m \leq n$. Then $\operatorname{rank}(\mathbf{X}) \leq r$ if and only if there exists a full row-rank matrix $\mathbf{R} \in \mathbb{R}^{(m-r) \times m}$ such that

$$
\begin{equation*}
\mathbf{R X}=\mathbf{0} \text {. } \tag{7.16}
\end{equation*}
$$

Let us suppose, without loss of generality, that $N_{e}+N_{B} \leq N$. Then, the theorem is applied for $\mathbf{X}=\left[\mathbf{B}\left|\boldsymbol{c}_{1}\right| \ldots \mid \boldsymbol{c}_{N_{c}}\right]^{T} \in \mathbb{R}^{\left(N_{c}+N_{B}\right) \times N}$ and $r=N_{B}$, thus $\mathbf{R} \in \mathbb{R}^{N_{c} \times\left(N_{c}+N_{B}\right)}$. Since the full row-rank condition still involves the rank, it is convenient to replace it with the orthogonality condition $\mathbf{R R}^{T}=\mathbf{I}_{N_{e}}$, where $\mathbf{I}_{N_{e}}$ denotes the $N_{e} \times N_{e}$ identity. It is evident that orthogonality implies full row-rank. Thus problem (7.15) can be recast as

$$
\begin{array}{ll}
\text { minimize } & \left\|\mathbf{R}\left[\mathbf{B}\left|\boldsymbol{c}_{1}\right| \ldots \mid \boldsymbol{c}_{N_{e}}\right]^{T}\right\|^{2} \\
\text { subject to } & \left\{\begin{array}{l}
\mathbf{R R}^{T}=\mathbf{I} \\
(\Delta \mathbf{M}, \Delta \mathbf{K}) \in \Gamma,
\end{array}\right. \tag{7.17}
\end{array}
$$

where $\|\cdot\|$ denotes the matrix 2 -norm. The resulting problem can be numerically solved with appropriate algorithms that are implemented in commercial software. Since the problem is nonconvex, the choice of the initial guess for the optimization variables is of fundamental importance. Two alternatives are proposed.

Initial guess 1. Concerning matrices $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$, a natural choice is to set both the modification matrices to be zero, in order to obtain structural modifications that are close to the original system design.

Concerning matrix $\mathbf{R}$, this work proposes a strategy that has been tested with favorable results. Let us denote $\boldsymbol{c}_{i}^{0}$ the vectors (7.12) where $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$ are set to zero, in accordance with the initial guess for the optimization variables. Let us define $\mathbf{X}^{0}:=\left[\mathbf{B}\left|\boldsymbol{c}_{1}^{0}\right| \ldots \mid \boldsymbol{c}_{N_{e}}^{0}\right]$ and compute the singular value decomposition of its transpose:

$$
\begin{equation*}
\left(\mathbf{X}^{0}\right)^{T}=\mathbf{U} \Sigma \mathbf{V}^{T} . \tag{7.18}
\end{equation*}
$$

The smallest $N_{e}$ singular values are approximated to zero, leading to a new matrix $\tilde{\boldsymbol{\Sigma}}$ whose rank is $N_{B}$. Thus a rank- $N_{B}$ approximation of $\mathbf{X}^{0}$ can be defined as

$$
\begin{equation*}
\left(\tilde{\mathbf{X}}^{0}\right)^{T}=\mathbf{U} \tilde{\mathbf{\Sigma}} \mathbf{}^{T} \tag{7.19}
\end{equation*}
$$

Finally, the initial guess for $\mathbf{R}$, denoted $\mathbf{R}^{0}$, is chosen in such a way that the column of $\left(\mathbf{R}^{0}\right)^{T}$ are an orthogonal basis of $\operatorname{Ker}\left(\left(\tilde{\mathbf{X}}^{0}\right)^{T}\right)$. Such a heuristic choice is justified by the fact that $\mathbf{R}^{0} \tilde{\mathbf{X}}^{0}=\mathbf{0}$, and since it is fair to assume $\mathbf{R}^{0} \tilde{\mathbf{X}}^{0} \approx \mathbf{R}^{0} \mathbf{X}^{0}$, the cost function of (7.17) can be expected to be small since the first iteration.

Initial guess 2. As a second choice, matrices $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$ can be chosen as the ones that allows the assignment of $\left(\omega_{i}, \boldsymbol{u}_{i}\right)$, for $i=1, \ldots, N_{e}$, by means of DSM, computed using the convex optimization problem proposed in [23]. As far as $\mathbf{R}$ is concerned, instead, the same strategy proposed in the previous approach should be employed, provided that the vectors $\boldsymbol{c}_{i}^{0}$ are defined accordingly to the choice of $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$.

## Algorithm 2

The second algorithm exploits the following theorem, which is known as the semidefinite embedding lemma [107].

Theorem 2. Let $\mathbf{X} \in \mathbb{R}^{m \times n}$, with arbitrary $m$ and $n$. Then $\operatorname{rank}(\mathbf{X}) \leq r$ if and only if there exist matrices $\mathbf{Y}=\mathbf{Y}^{T} \in \mathbb{C}^{m \times m}$ and $\mathbf{Z}=\mathbf{Z}^{T} \in \mathbb{C}^{n \times n}$ such that

$$
\operatorname{rank}(\mathbf{Y})+\operatorname{rank}(\mathbf{Z}) \leq 2 r, \quad\left[\begin{array}{cc}
\mathbf{Y} & \mathbf{X}  \tag{7.20}\\
\mathbf{X}^{T} & \mathbf{Z}
\end{array}\right] \succeq 0
$$

where the symbol $\succeq$ means that the matrix must be positive semidefinite.

Such a theorem allows to recast the rank minimization of a general rectangular matrix to the rank minimization of a square symmetric matrix, i.e. $\operatorname{diag}(\mathbf{Y}, \mathbf{Z})$. In such a way, it can be employed the popular trace heuristic [108], which consists in replacing $\operatorname{rank}(\operatorname{diag}(\mathbf{Y}, \mathbf{Z}))$ with $\operatorname{trace}(\operatorname{diag}(\mathbf{Y}, \mathbf{Z}))$ in the objective function of the optimization problem. An intuitive explanation of the trace heuristic is that a small trace forces the numerical rank to be low in the case of symmetric positive semidefinite matrices. In fact, for such matrices eigenvalues are real
and greater than, or equal to zero. Hence a small trace (i.e. the sum of the eigenvalues) is usually associated to a low rank (i.e. few non-zero eigenvalues). A rigorous interpretation of the trace heuristic can be found in [109].

Using Theorem 2 for $\mathbf{X}=\left[\mathbf{B}\left|\boldsymbol{c}_{1}\right| \ldots \mid \boldsymbol{c}_{N_{e}}\right]$ and $r=N_{B}$, the Eigenstructure Assignment problem (7.15) can be therefore recast as follows:

$$
\begin{array}{ll}
\operatorname{minimize} & \operatorname{trace} \operatorname{diag}(\mathbf{Y}, \mathbf{Z}) \\
\text { subject to } & \left\{\begin{array}{cc}
\mathbf{Y} & {\left[\mathbf{B}\left|\boldsymbol{c}_{1}\right| \ldots \mid \boldsymbol{c}_{N_{e}}\right]} \\
{\left[\mathbf{B}\left|\boldsymbol{c}_{1}\right| \ldots \mid \boldsymbol{c}_{N_{e}}\right]^{T}} & \mathbf{Z}
\end{array}\right] \succeq 0  \tag{7.21}\\
(\Delta \mathbf{M}, \Delta \mathbf{K}) \in \Gamma .
\end{array}
$$

The obtained optimization problem is convex if the feasibility constraint set $\Gamma$ is chosen to be a convex set, as usually done to represent this type of constraint [23]. Indeed, semidefinite constraints are convex and can be handled with several efficient and reliable software packages.

It should be pointed out that the proposed method, independently of the used algorithm, does not assure that the allowable subspaces of the modified system will actually contain the desired eigenvectors. In fact, the presence of the constraints expressed by $\Gamma$ and the dimension of $\mathbf{B}$ often prevent the exact attainment of the desired objective. In such cases, it is still necessary to assign by active control the projection of the desired eigenvector onto the allowable subspace of the modified system, which is denoted $\hat{\Psi}$. Such a projection is denoted $\boldsymbol{u}_{i, p}$ and is computed as

$$
\begin{equation*}
\boldsymbol{u}_{i, p}=\mathbf{L}_{i}\left(\mathbf{L}_{i}^{T} \mathbf{L}_{i}\right)^{-1} \mathbf{L}_{i}^{T} \boldsymbol{u}_{i} \tag{7.22}
\end{equation*}
$$

where $\mathbf{L}_{i}=\hat{\mathbf{H}}\left(\omega_{i}\right) \mathbf{B}$. However, the projection of $\boldsymbol{u}_{i}$ onto $\hat{\Psi}\left(\omega_{i}\right)$ approximates $\boldsymbol{u}_{i}$ in a better way than the projection onto $\Psi\left(\omega_{i}\right)$, leading to a more accurate attainment of the design specifications.

## Summary of method application

In conclusion, the proposed method can be summarized as follows. Given a set of desired eigenpairs $\left(\omega_{1}, \boldsymbol{u}_{1}\right), \ldots,\left(\omega_{N_{e}}, \boldsymbol{u}_{N_{e}}\right)$, the suitable modifications of the mass and stiffness matrices, i.e. $\Delta \mathbf{M}$ and $\Delta \mathbf{K}$, are sought. It is wanted that the allowable subspace contains the associated eigenvectors, in the modified system design. Namely,

$$
\begin{equation*}
\boldsymbol{u}_{i} \in \hat{\Psi}\left(\omega_{i}\right), \text { for all } i=1, \ldots, N_{e} \tag{7.23}
\end{equation*}
$$

The modifications are computed either solving the optimization problem (7.17) or (7.21), and in both cases the solution is constrained into a set $\Gamma$.

The desired eigenvectors are then projected as in Eq. (7.22) onto their respective allowable subspaces, and the control gain matrix $\mathbf{G}$ that allows to attain the desired closed-loop eigenpairs is computed. In principle, any method for active eigenstructure assignment could be employed, thus the appropriate method can be chosen according to specific design needs. For example, in this work it is used the method proposed in [35], which is capable of partial eigenstructure assignment, thus only $N_{e}$ selected open-loop eigenpairs are replaced with the desired ones, while the others are left unchanged, in such a way that spill-over is prevented. Moreover, such a method relies on a second-order system model, which is more suitable to vibration control with respect to the more common first-order model.

### 7.4. Numerical examples

### 7.4.1. Five degrees of freedom system

The first example is aimed at assessing the effectiveness of the proposed method and it will be proved that it can be used to achieve superior performances with respect to either sole passive control or sole active control. This example also justifies the need of a method specifically designed for hybrid control, performing concurrently synthesis of controller gains and structural modification, which offers better performances compared to the sequential combination of passive and active control.

A lumped parameter 5-degrees of freedom system is considered. This system has been used in several papers to validate some methods of Dynamic Structural Modification [9, 23, 44], and hence it is an interesting benchmark. Such a system, pictured in Figure 7.1, consists of 5 lumped masses. Each mass is connected to the adjacent ones and to the frame by linear springs. The original values of the masses and springs are shown in the first row of Table 7.1.


Figure 7.1. Five degrees of freedom system

Table 7.1. Mass and stiffness parameters of the five degrees of freedom system.

|  | $m_{1}[\mathrm{~kg}]$ | $m_{2}[\mathrm{~kg}]$ | $m_{3}[\mathrm{~kg}]$ | $m_{4}[\mathrm{~kg}]$ | $m_{5}[\mathrm{~kg}]$ | $k_{12}$ <br> $[\mathrm{kN} / \mathrm{m}]$ | $k_{23}$ <br> $[\mathrm{kN} / \mathrm{m}]$ | $k_{34}$ <br> $[\mathrm{kN} / \mathrm{m}]$ | $k_{45}$ <br> $[\mathrm{kN} / \mathrm{m}]$ | $k_{g, i}$ <br> $[\mathrm{kN} / \mathrm{m}]$ |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| Nominal <br> value | 1.34 | 2.61 | 8.21 | 5.12 | 1.73 | 82.1 | 73.5 | 68.2 | 73.6 | 98.9 |
| Lower <br> bound | 0 | 0 | 0 | 0 | 0 | - | - | - | - | 0 |
| Upper <br> bound | 2 | 2 | 2 | 2 | 2 | - | - | - | - | 483 |

The original system eigenstructure is shown in Table 7.2 (the meaning of the eigenvector components can be inferred from Figure 7.1). It is wanted that the system features the two eigenpairs shown in Table 7.3, which are those required in [23]. In order to achieve this result, it is supposed that the five masses and the five ground springs can be modified, according to the constraints specified in the second and third rows of Table 7.1. Moreover, it is supposed that two independent control forces $v_{1}$ and $v_{2}$ are exerted on the first two masses on the left, as shown in Figure 7.1.

Therefore, the dimension of the allowable subspaces is $N_{B}=2$, thus it is unlikely that active control suffices to assign accurately the required eigenvectors stated in Table 7.3. To assess if the desired eigenvectors lie in their corresponding allowable subspace, it is evaluated the cosine of the angle between each desired eigenvector and its projection onto the allowable subspace. The mode shape at 55 Hz can be assigned through active control with a good approximation, since the computed cosine is 0.9988 . In contrast, the mode shape at 39 Hz cannot be assigned and its projection onto the allowable subspace leads to poor results, since the computed cosine is 0.0447 .

To prove the superiority of the hybrid approach, the previously described eigenstructure assignment task is performed by means of the proposed method, as well as with state-of-the-art methods for either passive control [8] or active control [35]. Moreover, it has been attempted to apply sequentially the passive and the active method in order to demonstrate that such an approach, which could be also considered a hybrid method, is less effective than the proposed one, which tackles the problem in a concurrent way. Indeed, the structural modification of the proposed method are not aimed at directly assigning the desired eigenstructure but rather at properly modifying the allowable subspace. The two algorithms developed have been

Table 7.2. Original system eigenstructure.

| Mode number, $i$ | 1 | 2 | 3 | 4 | 5 |
| :--- | :---: | :--- | ---: | ---: | ---: |
| $f_{i}[\mathrm{~Hz}]$ | 22.28 | 32.61 | 42.92 | 52.71 | 64.57 |
| $u_{i, 1}$ | 0.244 | -0.222 | 0.983 | -0.019 | 1.000 |
| $u_{i, 2}$ | 0.460 | -0.337 | 1.000 | -0.008 | -0.482 |
| $u_{i, 3}$ | 1.000 | -0.417 | -0.218 | 0.025 | 0.032 |
| $u_{i, 4}$ | 0.673 | 1.000 | 0.060 | -0.235 | -0.004 |
| $u_{i, 5}$ | 0.358 | 0.737 | 0.094 | 1.000 | 0.003 |

Table 7.3. Desired eigenpairs.

| Mode number, $i$ | 1 | 2 |
| :--- | ---: | ---: |
| $f_{i}[\mathrm{~Hz}]$ | 39.00 | 55.00 |
| $u_{i, 1}$ | 0.05 | 1.00 |
| $u_{i, 2}$ | 0.00 | 0.80 |
| $u_{i, 3}$ | 0.20 | -0.10 |
| $u_{i, 4}$ | -0.55 | 0.01 |
| $u_{i, 5}$ | 1.00 | 0.00 |

implemented in Matlab, by exploiting the YALMIP [47] modeling language for optimization problems. The numerical solvers employed are fmincon from the Optimization Toolbox by Mathworks for the non-convex optimization of Eq. (7.17), and SeDuMi [64] for the semidefinite optimization of Eq. (7.21).

Table 7.4 shows the modifications of the physical parameters for all the methods involving structural modification. The quantities evaluated to assess the system performances are shown in Table 7.5. The first one is the absolute error of the obtained eigenfrequencies $\omega_{i, \text { obt }}$ with respect to the desired ones $\omega_{i, \text { des }}$, which should approach 0 . The second one is the cosine of the angle between each desired eigenvector $\boldsymbol{u}_{i, \text { des }}$ and the obtained one $\boldsymbol{u}_{i, \text { obt }}$, which should approach 1. Finally, the 2 -norm of the two modification matrices and the gain matrix are

Table 7.4. System modifications.

|  | DSM [23] | Hybrid method (Algorithm 1) |  | Hybrid method <br> (Algorithm 2) |
| :---: | :---: | :---: | :---: | :---: |
|  |  | Initial guess 1 | Initial guess 2 |  |
| $\overline{m_{l}[\mathrm{~kg}]}$ | 0.878 | 1.857 | 2.000 | 0.000 |
| $m_{2}[\mathrm{~kg}]$ | 0.000 | 1.852 | 1.999 | 0.965 |
| $m_{3}[\mathrm{~kg}]$ | 0.808 | 0.045 | 0.239 | 0.000 |
| $m_{4}[\mathrm{~kg}]$ | 1.495 | 1.767 | 1.999 | 2.000 |
| $m_{5}[\mathrm{~kg}]$ | 1.816 | 1.908 | 1.998 | 1.927 |
| $k_{g, 1}[\mathrm{kN} / \mathrm{m}]$ | 149.69 | 28.948 | 44.360 | 47.103 |
| $k_{8,2}[\mathrm{kN} / \mathrm{m}]$ | 153.20 | 22.485 | 150.636 | 258.234 |
| $k_{\text {g,3 }}[\mathrm{kN} / \mathrm{m}]$ | 0.000 | 69.105 | 80.905 | 113.385 |
| $k_{8,4}[\mathrm{kN} / \mathrm{m}]$ | 11.94 | 14.094 | 28.096 | 22.897 |
| $k_{\text {g, }}[\mathrm{kN} / \mathrm{m}]$ | 0.000 | 5.421 | 10.861 | 4.355 |

evaluated to provide a concise measure of the amount of system modifications and of the magnitude of the controller gains.

Table 7.5 enables to evaluate the attainment of the specifications for each one of the considered methods for eigenstructure assignment. The method of [23], that exploits DSM, can accurately assign both the eigenvectors, but the assignment of eigenfrequencies is missed for both the modes, although the absolute error could be acceptable for some applications. In contrast, the other methods that involve active control are capable of an exact assignment of the eigenvalues, as expected being the system controllable. However, the method that exploits only feedback control can assign with satisfactory accuracy only the eigenvector of the 55 Hz mode. In contrast, the obtained eigenvector at 39 Hz is almost orthogonal to the desired one and hence the design specifications are missed.

Table 7.5. Eigenstructure comparison.

|  | $\begin{gathered} \text { DSM } \\ {[23]} \end{gathered}$ | Active control [35] | DSM+Active control | Hybrid method (Algorithm 1) |  | Hybrid method (Algorithm 2) |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  |  |  |  | Initial guess 1 | Initial guess 2 |  |
| $\left\|\omega_{1, \text { des }}-\omega_{1, \text { obt }}\right\|_{(H z)}$ | 0.068 | 0.000 | 0.000 | 0.000 | 0.000 | 0.000 |
| $\cos \left(\boldsymbol{u}_{1, \text { des }}, \boldsymbol{u}_{1, \text { obt }}\right)$ | 0.9964 | 0.0447 | 0.9854 | 1.0000 | 1.0000 | 0.9941 |
| $\left\|\omega_{2, \text { des }}-\omega_{2, \text { obt }}\right\|_{(H z)}$ | 0.020 | 0.000 | 0.000 | 0.000 | 0.000 | 0.000 |
| $\cos \left(\boldsymbol{u}_{2, \text { des }}, \boldsymbol{u}_{2, \text { obt }}\right)$ | 0.9997 | 0.9988 | 0.9998 | 1.0000 | 0.9999 | 1.0000 |
| $\\|\Delta \mathbf{M}\\|$ | 2.6375 | - | 2.6375 | 3.6931 | 4.0052 | 2.9403 |
| $\\|\Delta \mathbf{K}\\|$ | $2.1452 \cdot 10^{5}$ | - | $2.1452 \cdot 10^{5}$ | $7.9668 \cdot 10^{4}$ | $1.7920 \cdot 10^{5}$ | $2.8689 \cdot 10^{5}$ |
| $\\|\mathbf{G}\\|$ | - | $2.4946 \cdot 10^{6}$ | $2.7771 \cdot 10^{4}$ | $9.7315 \cdot 10^{6}$ | $1.2541 \cdot 10^{6}$ | $4.2038 \cdot 10^{4}$ |

The results obtained suggest that a hybrid approach can actually combine the advantages of passive and active control. In fact, both the algorithms for hybrid control enable an accurate assignment of both eigenvalues and eigenvectors, which are almost parallel to the desired ones.

Among the three approaches that combine DSM and active control, however, the two that exploits the proposed rank minimization formulation are the ones that allow to attain the best fulfillment of the design specifications.

### 7.4.2. Thirty degrees of freedom system

The second example is devoted to assess the performances of the method in the presence of large dimensional model and to evaluate the effect of the number of the control force (and hence the underactuation degree, i.e. the difference between the number of degrees of freedom and the number of control inputs). The vibrating system studied is obtained by expanding the system of the previous example through the addition of 25 masses. The resulting system has 30 degrees of freedom. The mass and the stiffness added are shown in Table 7.6. All the ground springs are, instead, supposed equal to $k_{g, i}=98.9 \mathrm{kN} / \mathrm{m}, i=1, \ldots, 30$, as in the previous example. The modifiable parameters assumed are the masses and the ground springs, and the constraints on the physical modifications are the same as in the previous example, i.e. $0 \leq \Delta m_{i} \leq 2 \mathrm{~kg}$ and $0 \leq \Delta k_{g, i} \leq 4.83 \cdot 10^{5} \mathrm{kN} / \mathrm{m}$, for $i=1, \ldots, 30$. It is assigned only one eigenpair, for simplicity of

Table 7.6. Mass and stiffness parameters of the thirty degrees of freedom system.

|  | $m_{6}[\mathrm{~kg}]$ | $m_{7}[\mathrm{~kg}]$ | $m_{8}[\mathrm{~kg}]$ | $m_{9}[\mathrm{~kg}]$ | $m_{10}[\mathrm{~kg}]$ | $m_{11}[\mathrm{~kg}]$ | $m_{12}[\mathrm{~kg}]$ | $m_{13}[\mathrm{~kg}]$ | $m_{14}[\mathrm{~kg}]$ | $m_{15}[\mathrm{~kg}]$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| Nominal <br> value | 3 | 1 | 4 | 1 | 5 | 9 | 2 | 6 | 5 | 3 |

representation. In particular, the desired frequency is 39 Hz and the associated eigenvector is shown in Table 7.7. Such numbers have been chosen as uniformly distributed random numbers in the interval $[-1 ; 1]$.

Table 7.7. Desired eigenvector for the thirty degrees of freedom system.

| $\mathrm{u}_{1,1}$ | $\mathrm{u}_{1,2}$ | $\mathrm{u}_{1,3}$ | $\mathrm{u}_{1,4}$ | $\mathrm{u}_{1,5}$ | $\mathrm{u}_{1,6}$ | $\mathrm{u}_{1,7}$ | $\mathrm{u}_{1,8}$ | $\mathrm{u}_{1,9}$ | $\mathrm{u}_{1,10}$ |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| -0.7867 | 0.9238 | -0.9907 | 0.5498 | 0.6346 | 0.7374 | -0.8311 | -0.2004 | -0.4803 | 0.6001 |
| $\mathrm{u}_{1,11}$ | $\mathrm{u}_{1,12}$ | $\mathrm{u}_{1,13}$ | $\mathrm{u}_{1,14}$ | $\mathrm{u}_{1,15}$ | $\mathrm{u}_{1,16}$ | $\mathrm{u}_{1,17}$ | $\mathrm{u}_{1,18}$ | $\mathrm{u}_{1,19}$ | $\mathrm{u}_{1,20}$ |
|  |  |  |  |  |  |  |  |  |  |
| -0.1372 | 0.8213 | -0.6363 | -0.4724 | -0.7089 | -0.7279 | 0.7386 | 0.1594 | 0.0997 | -0.7101 |
| $\mathrm{u}_{1,21}$ | $\mathrm{u}_{1,22}$ | $\mathrm{u}_{1,23}$ | $\mathrm{u}_{1,24}$ | $\mathrm{u}_{1,25}$ | $\mathrm{u}_{1,26}$ | $\mathrm{u}_{1,27}$ | $\mathrm{u}_{1,28}$ | $\mathrm{u}_{1,29}$ | $\mathrm{u}_{1,30}$ |
| 0.7061 | 0.2441 | -0.2981 | 0.0265 | -0.1964 | -0.8481 | -0.5202 | -0.7534 | -0.6322 | -0.5201 |

The proposed method has been tested for different choices of the actuation, by actuating the $N_{B}$ masses on the left, with $N_{B}$ ranging from 1 to 30 . The computation of the structural modifications has been performed with both the Algorithms described in Section 7.3, choosing the Initial guess 2 for Algorithm 1. The performances are evaluated through the cosine of the angle between the desired eigenvector and its projection onto the modified allowable subspace $\hat{\Psi}$ and through the CPU time required to solve the problems. Such quantities are represented in Figure 7.2, as functions of the number of the actuation forces $N_{B}$.

The results obtained show that Algorithm 2 enables to shape more effectively the allowable subspace $\Psi$ and hence to approximate more effectively the requirements. In fact, the subspace modification computed by Algorithm 2 allows to assign an eigenvector that matches more closely the desired one, as proved by the higher cosine. For what concerns the CPU time, instead, the behaviour of the two algorithms is very different. The representation of the time versus $N_{B}$ suggests a strong correlation in the case of Algorithm 2, which is faster in the case of highly underactuated systems (i.e. those involving smaller matrices). In contrast, the CPU time for Algorithm 1 is less sensitive to the number of actuation forces, but it is certainly much faster when the system is mildly underactuated.


Figure 7.2. Comparison of the algorithms.
For the above reasons, Algorithm 2 should be preferred to achieve the best performance, unless a system with low underactuation degree is considered and CPU time is a critical issue.

### 7.5. Conclusion

In this Chapter, a limitation of active vibration control that affects the capability of assigning eigenvectors is addressed. It is demonstrated that, with a proper adjustment of the inertial and/or elastic parameters of the system, such a limitation can be overcome. It is proposed a method to numerically compute the optimal values of such parameters, that consists in the minimization of the rank of a matrix. Two algorithms have been developed to solve such a problem, each one relying on a different mathematical framework.

The proposed approach has been numerically validated in two cases of interest. The first testcase aims to assign two eigenpairs to a five degrees of freedom system. Such an objective can be fairly accomplished with techniques of Dynamic Structural Modification available in the literature. It is nevertheless demonstrated that the proposed hybrid approach outperforms both the methods of DSM and those of active control. Moreover, it is shown that sequentially employing DSM and active control actually results in a good eigenstructure assignment, but it is still inferior with respect to the proposed concurrent method. The second test-case concerns
a single-mode assignment to a 30 degrees of freedom system, in different actuation conditions. In such a way, it has been possible also to compare the two developed algorithms and to validate the method in a large-scale example.

## Chapter 8. Experimental application of Hybrid Control

In this Chapter it is presented an experimental application of the hybrid control technique introduced in the thesis. In particular, it is considered a cantilever beam which is excited by a shaker that acts as a source of disturbance. It is wanted to confine the vibrations to the part of the beam near the free end. A piezoelectric actuator is used for this purpose, but it is proved that only one control force is not sufficient to achieve a satisfactory result. Exploiting the hybrid technique of the previous Chapter, the closed-loop eigenstructure is significantly improved.

### 8.1. Description of the system

The system investigated is a cantilever beam clamped on a slip-table which is actuated by an electrodynamic shaker. The geometric and physical properties of the beam are shown in Table 8.1. The beam is modelled through 8 Euler-Bernoulli beam elements, whose lengths are shown in Figure 8.1. The resulting system model, supposed that the slip-table cannot move, has $N=16$


Figure 8.1. Model of the cantilever beam.
Table 8.1. Nominal values of the physical parameters of the beam.

| Parameter | Value |
| :--- | :--- |
| Beam length $L[\mathrm{~mm}]$ | 495 |
| Beam thickness $h[\mathrm{~mm}]$ | 4 |
| Beam width $b[\mathrm{~mm}]$ | 40 |
| Young's modulus $E[\mathrm{GPa}]$ | 55 |
| Density $\rho\left[\mathrm{kg} / \mathrm{m}^{3}\right]$ | 2400 |

degrees of freedom. It is considered that at selected locations of the beam can be placed two masses of the desired magnitude.

### 8.2. The piezoelectric actuator

A piezoelectric actuator is employed for vibration control. The actuator has been placed on the desired location of the beam by bonding it to the beam surface. Linear theory has been adopted to model it, even if some nonlienarities such as hysteresis and variable gains might occur, and it is assumed that there is no slip at the interface. The fundamental properties of the actuator are shown in Table 8.2. Consistently with the literature on piezoelectric actuators [110, 111], the assumption that the applied voltage is uniform over the element allows to represent the force exerted by the control device as two opposite torques of magnitude $F_{C}$, as represented in Figure 8.2.

Table 8.2. Nominal value of the physical parameters of the piezoelectric actuator.

| Parameter | Value |
| :--- | :--- |
| Length $L_{p}[\mathrm{~mm}]$ | 61 |
| Thickness $h_{p}[\mathrm{~mm}]$ | 0.8 |
| Beam width $b_{p}[\mathrm{~mm}]$ | 35 |
| Young's modulus $\mathrm{E}_{\mathrm{p}}[\mathrm{GPa}]$ | 34.7 |
| Mass $M_{p}[\mathrm{~g}]$ | 7.2 |



Figure 8.2. Passive and active modification of the beam.

These torques are, in turn, proportional to the applied voltage, with a gain that is approximately constant and is equal to $0.034 \mathrm{Nm} / \mathrm{V}$. In practice, the experimental characterization of the actuator has shown that this gain might oscillate about this nominal value (the highest deviation
is about the 20\%), because of the aforementioned nonlinear behaviours. Finally, the contribution of the actuator to the finite element model has been represented through additive mass and stiffness matrices based on the Euler-Bernoulli linear theory.

### 8.3. The system model

The resulting system model can be represented by suitable mass matrix $\mathbf{M} \in \mathbb{R}^{16 \times 16}$ and stiffness matrix $\mathbf{K} \in \mathbb{R}^{16 \times 16}$ whose entries can be determined exploiting the classical theories of EulerBernoulli beam and Finite Element modelling, by also correcting the components that are affected by the local increase of inertia and stiffening due to the presence of the actuator and of the additional masses.

Damping can be taken into account employing the proportional Rayleigh damping, namely it is defined

$$
\begin{equation*}
\mathbf{C}=\alpha \mathbf{M}+\beta \mathbf{K} \tag{8.1}
\end{equation*}
$$

for the desired values of $\alpha$ and $\beta$. In this work it is chosen $\alpha=1.05$ and $\beta=25 \cdot 10^{-6}$. The Rayleigh coefficients have been estimated by identifying the damping through the half power method on the FRF relating the actuator torque and the strain gauge response.

For simplicity, the coordinates of the system are collected in a vector

$$
\begin{equation*}
\boldsymbol{u}=\left\{z_{1}, \vartheta_{1}, z_{2}, \vartheta_{2}, \ldots, z_{8}, \vartheta_{8}\right\}^{T} \tag{8.2}
\end{equation*}
$$

with the meaning of the symbols that can be inferred from Figure 8.1. The case in which the motion of the slip-table is allowed is represented in Figure 8.3.


Figure 8.3. Model of the moving beam.

The overall motion is thus represented by the augmented system of $N+1$ differential equation

$$
\left[\begin{array}{cc}
\mathbf{M} & \mathbf{M S}  \tag{8.3}\\
\mathbf{S}^{T} \mathbf{M} & \mathbf{S}^{T} \mathbf{M S}+M_{C}
\end{array}\right]\left\{\begin{array}{c}
\ddot{\boldsymbol{u}} \\
\ddot{\boldsymbol{x}}
\end{array}\right\}+\left[\begin{array}{cc}
\mathbf{C} & \mathbf{0} \\
\mathbf{0} & 0
\end{array}\right]\left\{\begin{array}{c}
\dot{\boldsymbol{u}} \\
\dot{\boldsymbol{x}}
\end{array}\right\}+\left[\begin{array}{cc}
\mathbf{K} & \mathbf{0} \\
\mathbf{0} & 0
\end{array}\right]\left\{\begin{array}{c}
\boldsymbol{u} \\
\boldsymbol{x}
\end{array}\right\}=\left\{\begin{array}{l}
\mathbf{0} \\
1
\end{array}\right\} F_{A}+\boldsymbol{b} F_{C}
$$

where $F_{A}$ is the force exerted by the electrodynamic shaker on the slip-table (whose mass is $M_{C}$ ), the coordinate $x$ represents the displacement of the slip-table and $\mathbf{S}$ is the matrix of the nodal sensitivity coefficients with respect to $x$ :

$$
\begin{equation*}
[\mathbf{S}]_{i}=\frac{\partial u_{i}}{\partial x}, \quad i=1, \ldots, 2 N \tag{8.4}
\end{equation*}
$$

Since $\frac{\partial z_{i}}{\partial x}=1$ and $\frac{\partial \vartheta_{i}}{\partial x}=0$, for $i=1, \ldots, N$, then $\mathbf{S}=\{1,0,1, \ldots\}^{T}$. It should be noted that the model obtained recalls the typical ERLS based model, already introduced in Chapter 5 to represent flexible mechanisms. In this case the rigid body motion is represented by the translation of the slip table, which also sets the ERLS.

Since $M_{C} \gg \mathbf{S}^{T} \mathbf{M S}$ and being the slip-table controlled through a closed-loop control scheme with high gains, then the beam has neglibible effect on the slip-table. Hence it is convenient to consider the $N$-dimensional model only, with $x$ as an independent input. Hence,

$$
\begin{equation*}
\mathbf{M} \ddot{\boldsymbol{u}}+\mathbf{C} \dot{\boldsymbol{u}}+\mathbf{K} \boldsymbol{u}=\boldsymbol{b} F_{C}-\mathbf{M S} \ddot{x} . \tag{8.5}
\end{equation*}
$$

Finally, the model can be recast in the first-order form that is suitable for the observer design:

$$
\left\{\begin{array}{c}
\ddot{\boldsymbol{u}}  \tag{8.6}\\
\dot{\boldsymbol{u}}
\end{array}\right\}=\left[\begin{array}{cc}
-\mathbf{M}^{-1} \mathbf{C} & -\mathbf{M}^{-1} \mathbf{K} \\
\mathbf{I} & \mathbf{0}
\end{array}\right]\left\{\begin{array}{c}
\dot{\boldsymbol{u}} \\
\boldsymbol{u}
\end{array}\right\}+\left\{\begin{array}{c}
\mathbf{M}^{-1} \boldsymbol{b} \\
\mathbf{0}
\end{array}\right\} F_{C}+\left\{\begin{array}{c}
-\mathbf{S} \\
\mathbf{0}
\end{array}\right\} \ddot{\boldsymbol{x}} .
$$

### 8.4. Hybrid control

The six lowest frequencies of the system are shown in Figure 8.4. It is wanted to control the beam in such a way that it features a mode at 50 Hz whose associated mode shape is pictured in Figure 8.5, together with the closest mode of the original system design, namely the second lowest frequency one. It is wanted to employ the piezoelectric actuator to exert the suitable control forces to obtain the desired mode shape. Trying all the possible configurations, it has been concluded that, in order to improve controllability, it is recommendable to place the actuator in correspondence of the sixth Finite Element from the left, as in Figure 8.2.


Figure 8.4. Lowest frequency modes of the beam.


Figure 8.5. Original ( 79.95 Hz ) and desired $(50 \mathrm{~Hz})$ modes.

Nevertheless, although the closed-loop pole at 50 Hz can be actually attained, it is not possible to achieve the desired mode shape, because the corresponding eigenvector does not belong to the allowable subspace. The computed cosine of the angle between the obtained eigenvector and the desired one, in fact, is 0.3767 , which is very far from the target value 1 . In Figure 8.6 are represented the two eigenvectors and it is clear that the results obtained by state feedback control are not satisfactory.


Figure 8.6. Mode attainable by active control.
Given the possibility to add two masses on the beam, it is used the hybrid control method described in the previous Chapter to suitably modify the allowable subspace. The values of the computed masses are shown in Table 8.3. The modified system can be controlled more easily: the closed-loop pole at 50 Hz can be obtained and the best assignable eigenvector is almost parallel to the desired one, being the cosine of the angle between them equal to 0.9929 . Such an accurate match can be also assessed through the representation of the obtained mode shape in Figure 8.7.

Table 8.3. Structural modifications.

| Mass | Constraints | Optimal value |
| :--- | :--- | :---: |
| $m_{1}[\mathrm{~g}]$ | $[0,200]$ | 70.6 |
| $m_{2}[\mathrm{~g}]$ | $[0,200]$ | 200.0 |

The masses made of steel have been made to approximate the optimal computed masses. Since concentrated masses do not provide an accurate modelling of the resulting system, it has been


Figure 8.7. Mode attainable by hybrid control.
chosen to consider the masses as rigid bodies with a rotational inertia. The mass of $m_{1}$, inclusive of the screws, is 63.9 g , and the rotational inertia is $J_{1}=1.185 \cdot 10^{-5} \mathrm{~kg} \cdot \mathrm{~m}^{2}$. Analogously, $m_{2}=202.9 \mathrm{~g}$ and $J_{2}=6.964 \cdot 10^{-5} \mathrm{~kg} \cdot \mathrm{~m}^{2}$. With these modifications, the best assignable eigenvector does not change appreciably.

### 8.5. State estimation

In order to employ a state-feedback control, it is necessary to develop a state estimator that provides an effective approximation of the unmeasurable state from a reduced set of appropriate measurements. In order to ensure adequate observability and hence to grant the existence of reliable estimates even in the presence of modelling errors and measurement noises, the two following measurements have been chosen as the sensed output:

- a pair of resistive strain gauges (in a half-bridge configuration), which provide the local strain, and hence the curvature $\mathcal{E}$ of the beam. Curvature is defined through the shape function and is a linear combination of the displacements of the nodes of the beam:

$$
\varepsilon=\left\{-\frac{6}{l^{2}}+\frac{12 s}{l^{3}}-\frac{4}{l}+\frac{6 s}{l^{2}} \quad \frac{6}{l^{2}}-\frac{12 s}{l^{3}}-\frac{2}{l}+\frac{6 s}{l^{2}}\right\}\left\{\begin{array}{c}
z_{i}  \tag{8.7}\\
\vartheta_{i} \\
z_{i+1} \\
\vartheta_{i+1}
\end{array}\right\}
$$

where the meaning of the symbols can be found in Figure 8.8.


Figure 8.8. Definition of curvature

- a laser vibrometer employed to measure directly the velocity of a node.

The position of the two sensors has been chosen to maximize observability and to improve the signal-to-noise ratio. In particular, in Figure 8.9 are plotted the values of the curvature in the middle of each finite element, for both the lowest frequency open-loop mode (on the left) and the best achievable mode at 50 Hz (on the right). It has been chosen to position the strain gauge in the middle of the fourth Finite Element, namely at roughly 175 mm from the free end of the beam. The laser vibrometer, instead, has been positioned to measure the velocity of a point near the free end of the beam, which is the part of the system that vibrates the most.


Figure 8.9. Curvatures of the two modes of interest

The estimator that has been implemented is a linear Kalman filter based on a reduced-order model on a modal base. Such a choice is motivated by the need of reducing the computational effort for the real-time solution of the observer differential equations. Indeed, adopting the fullorder model with 32 dofs would have required a relevant CPU time for each time step and, at the same time, a reduced time step to avoid aliasing of the high frequency modes and the numerical instability due to the presence of such high frequency dynamics. On the other hand, high frequency modes are weakly observable, given the low-pass filters on the sensors and given the presence of high-frequency noise, and weakly excited in the range of frequency of interest for the application developed.

The first-order formulation of the system model can be written in the following compact form by defining the nodal state vector as $\chi=\left\{\begin{array}{ll}\dot{\boldsymbol{u}}^{T} & \boldsymbol{u}^{T}\end{array}\right\}^{T}$ :

$$
\left\{\begin{array}{l}
\dot{\chi}=\mathbf{A} \chi+\mathbf{B} u  \tag{8.8}\\
y=\mathbf{C} \chi
\end{array}\right.
$$

The model in (8.8) is recast in the modal canonical form by using the linear transformation $z=\mathbf{T} \boldsymbol{\chi}\left(\mathbf{T} \in \mathbb{R}^{2 N \times 2 N}\right)$ where $z$ are the modal coordinates of the first-order model :

$$
\left\{\begin{array}{l}
\dot{z}=\mathbf{A}_{M} z+\mathbf{B} u  \tag{8.9}\\
y=\mathbf{C}_{M} z
\end{array}\right.
$$

The matrices of the transformed system are $\mathbf{A}_{M}=\mathbf{T} \mathbf{A T}^{-1}, \mathbf{B}_{M}=\mathbf{T B}$ and $\mathbf{C}_{M}=\mathbf{C T}^{-1}$. In particular, $\mathbf{A}_{M}$ has $2 \times 2$ blocks on the diagonal for each vibrational mode and it is zero everywhere else. Thus there is a block $\left[\begin{array}{cc}\sigma & v \\ -v & \sigma\end{array}\right]$ for any complex conjugate pair of eigenvalues $\sigma \pm i v$.

The reduction consists in subdividing the modal base $z$ into a set of neglected modes $z_{N}$, that will be discarded in the model adopted for the observer, and a small set of retained ones $z_{R}$, i.e.

$$
\left\{\begin{array}{l}
\left\{\begin{array}{l}
\dot{z}_{N} \\
\dot{z}_{R}
\end{array}\right\}=\left[\begin{array}{cc}
\mathbf{A}_{N} & \mathbf{0} \\
\mathbf{0} & \mathbf{A}_{R}
\end{array}\right]\left\{\begin{array}{l}
z_{N} \\
z_{R}
\end{array}\right\}+\left\{\begin{array}{l}
\mathbf{B}_{N} \\
\mathbf{B}_{R}
\end{array}\right\} u  \tag{8.10}\\
y=\left[\begin{array}{ll}
\mathbf{C}_{N} & \mathbf{C}_{R}
\end{array}\right]\left\{\begin{array}{l}
z_{N} \\
z_{R}
\end{array}\right\}
\end{array}\right.
$$

The block-diagonal form of the model allows the straightforward model reduction by truncating the neglected modes, leading to the reduced order model

$$
\left\{\begin{array}{l}
\dot{z}_{R}=\mathbf{A}_{R} z_{R}+\mathbf{B}_{R} u  \tag{8.11}\\
y=\mathbf{C}_{R} z_{R} .
\end{array}\right.
$$

The state estimator is based on the classical Luenberger closed-loop observer on a modal base:

$$
\left\{\begin{array}{l}
\hat{\dot{z}}_{R}=\mathbf{A}_{R} \hat{z}_{R}+\mathbf{B}_{R} u+\mathbf{L}(y-\hat{y})  \tag{8.12}\\
\hat{y}=\mathbf{C}_{R} \hat{z}_{R}
\end{array}\right.
$$

with $\mathbf{L}$ being the estimator gain and the hat ${ }^{\wedge}$ marking the estimated quantities. In the present work, it has been chosen to retain just the 3 lowest frequency modes. The neglected state variables are instead estimated as zero: $\hat{z}_{N}=\mathbf{0}$ and $\hat{\dot{z}}_{N}=\mathbf{0}$

Finally, the estimated state on the nodal base, as required by the controller, is computed through the inverse of the transformation matrix

$$
\hat{\chi}=\mathbf{T}^{-1} \hat{z}=\mathbf{T}^{-1}\left\{\begin{array}{l}
\hat{z}_{N}  \tag{8.13}\\
\hat{z}_{R}
\end{array}\right\}=\mathbf{T}^{-1}\left\{\begin{array}{c}
\mathbf{0} \\
\hat{z}_{R}
\end{array}\right\}
$$

where the neglected modes have been considered equal to zero.
The observer gains have been computed through the Kalman approach, so as to minimize the mean square error between the estimated and the actual values of the state variables. Since the problem is linear (except for the presence of some nonlinearities due to the actuators), a stationary Kalman filter has been adopted and a time-infinite Riccati equation has been solved to compute the gain starting from an estimation of the measurement and process noise covariance matrices.

### 8.6. Model identification and updating

In order for the estimator to reliably approximate the state of the system, it is important that the model of the system is accurate. Therefore, the system parameters have been tuned in such a way that the model features a set of natural frequencies that match the measured ones, shown in Figure 8.10.

It has been chosen to use the method of Chapter 4 to perform model updating, in such a way that only the lowest frequencies are adjusted without affecting the highest. Moreover, such a method enables to modify only the parameters of interest, thus assuring that the updated model is physically meaningful. In particular, it has been chosen to update the quantities that are more prone to measurement errors, namely the Young's modulus of the beam $E$, the density of the beam $\rho$ and the rotational inertias $J_{1}$ and $J_{2}$.


Figure 8.10. Measured natural frequencies.

The computed update parameters are shown in Table 8.4. The obtained natural frequencies can be compared with the original and the measured ones in Table 8.5 , proving that the updated model provides a better representation of the actual system. The accordance of the model with the measured data can be assessed also in Figure 8.11 which compares the forced response of the numerical model and of the real system through the FRFs whose input is the control voltage and the output is the strain measured by the strain gauge. It is evident the capability of the model to provide a very precise representation of the two lowest frequency modes and of the system dynamic in the range of frequencies shown.

Table 8.4. Model updating of the system.

| Parameter | Original value | Updated value |
| :--- | :--- | :--- |
| Young's modulus E [GPa] | 55 | 64 |
| Density $\rho\left[\mathrm{kg} / \mathrm{m}^{3}\right]$ | 2400 | 2450 |
| Rotational inertia $J_{l}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | $1.185 \cdot 10^{-5}$ | $1.185 \cdot 10^{-5}$ |
| Rotational inertia $J_{2}\left[\mathrm{~kg} \cdot \mathrm{~m}^{2}\right]$ | $6.964 \cdot 10^{-5}$ | $6.964 \cdot 10^{-5}$ |

Table 8.5. Frequency comparison.

|  | Original <br> model | Experimental | Updated <br> model |
| :--- | ---: | ---: | ---: |
| $f_{1}[\mathrm{~Hz}]$ | 7.291 | 7.7 | 7.823 |
| $f_{2}[\mathrm{~Hz}]$ | 38.887 | 42.0 | 41.795 |
| $f_{3}[\mathrm{~Hz}]$ | 172.336 | 169.7 | 183.816 |
| $f_{4}[\mathrm{~Hz}]$ | 295.626 | 307.1 | 316.296 |
| $f_{5}[\mathrm{~Hz}]$ | 477.269 | 487.5 | 512.071 |



Figure 8.11. Comparison of the measured receptances with the ones of the updated model Once that the model has been updated, the hybrid control synthesis has been corrected to account for the actual model parameters that slightly differ from the nominal one. The best assignable eigenvector of the updated system at 50 Hz is plotted in Figure 8.12, in comparison with the desired one. Despite the difference between the two eigenvectors being appreciable, the attainable one is valid for the purpose of vibration confinement, as proved by the computed cosine, which is 0.9917 .


Figure 8.12. Mode attainable by hybrid control after model updating

### 8.7. Experimental setup

The control scheme is implemented on a PC where a real-time kernel that interfaces with the operating system is installed through The MathWorks Real-Time Windows Target. In order to trade between high-frequency control and low computational effort, the sampling time adopted is 1 ms , and the numerical discretization of the differential equations governing the state observer has been performed through the $4^{\text {th }}$ order Runge-Kutta scheme.

The estimator has been tuned so that to ensure precise and undelayed estimates while filtering noise. An example of the result obtained is shown in Figure 8.13, which compares the estimated and the actual values of the strain gauge and of the vibrometer and proves the correctness of the estimation strategy proposed in this work.


Figure 8.13. Comparison of actual and estimated quantities

The devices that constitutes the experimental setup are shown in the following.

- E-413.D2 Power amplifier for piezoelectric patch transducers (Figure 8.14);
- Input voltage range:-2 to 8 V
- Output voltage range: -100 to 400 V


Figure 8.14. Power amplifier.

- Brüel \& Kjær NEXUS 2693, low noise conditioning amplifier for accelerometers (Figure 8.15). The signals have been filtered through a low-pass second order filter with 1000 Hz bandwidth.


Figure 8.15. Conditioning amplifier for accelerometers.
The following accelerometers have been processed (see Section 5.5. for more details on the sensors):

- Endevco 27AM1-100 10203, miniaturized ICP accelerometer, adopted to verify local acelerations or speed (through integration performed by the Nexus amplifier)
- Brüel \& Kjær 4508 ICP accelerometer, to measure the acceleration of the sliptable to be input to the observer
- HBM QuantumX MX878B + QuantumX MX840B, multichannel data acquisition system for conditioning the strain gauge bridge (Figure 8.16). To eliminate measurement noise and to prevent aliasing problems, the strain-gauge signals are passed through a low-pass filter cut-off frequency of 500 Hz .


Figure 8.16. Data acquisition system.

- Polytec OFV 525 laser vibrometer (Figure 8.17) with OFV 5000 controller, to measure the nodal speed employed for the synthesis of the state observer


Figure 8.17. Laser vibrometer.

- National Instruments CB-68LP: unshielded Input/Output accessories with 68 screw terminals (Figure 8.18) for easy signal connection to the multifunction Input/Output board 6070 E installed on the PXI adopted, which implements the control scheme running in Matlab-Simulink


Figure 8.18. Unshielded 68-Pin I/O connector block.

- Tira tbs 6000 permanent magnet shaker, to provide external base excitation.
- Dytran 3136A1945 ICP accelerometer, employed for the closed loop control of the shaker.
- Frontend system LMS SCADAS Recorder (Figure 8.19), running the software LMS Test.Lab Sine Vibrational Control to drive the shaker.


Figure 8.19. Drive of the shaker and frontend system

In Figure 8.20 the instrumented beam is shown.


Figure 8.20. Experimental setup.

### 8.8. Experimental results

The obtained experimental results show that the employed control scheme actually succeeded in obtaining the desired closed-loop pole at the frequency of 50 Hz , as shown in Figure 8.21, where frequency response functions of the system with control and without control are plotted in the frequency range of interest.


Figure 8.21. Comparison of the FRF with control and without control.
Also the desired mode shape have been achieved with a good accuracy, as it can be assessed from Figure 8.22 and 8.23 , where the obtained mode shape is compared with, respectively, the desired one and the best theoretically attainable mode shape. The cosine of the angle between


Figure 8.22. Comparison of the desired mode with the experimentally obtained mode.


Figure 8.23. Comparison of the achievable mode with the experimentally obtained mode.
the obtained mode and the desired one is 0.9315 , while the cosine of the angle between the obtained mode and the best achievable eigenvector is 0.9906 .

In order to improve the results, it would certainly be beneficial to employ a more accurate model of the system. In fact, it has been possible to adequately characterize the lowest frequency modes, especially after the updating procedure, but it would be desirable to reduce even further the uncertainty that unavoidably affects the model. In particular, the piezoelectric effect of the actuator has been only approximatedly described, since it has been observed a mildly non-linear behaviour that has been neglected, for simplicity. Other aspects of the model that can be improved are the local stiffening due to the presence of the actuator and the representation of the masses $m_{1}$ and $m_{2}$, that are supposed for simplicity to be lumped masses at the design stage, while the rotational inertia is considered only afterwards. Moreover, the state observer, although it has been proved to provide a reliable estimation of the state, could still benefit from a more accurate system model. These and other improvements may help to achieve experimental results that are closer to the numerical prediction, but even the results in the present Section can be
considered very satisfactory, since they enable to assess that the developed hybrid technique can be successfully applied to the control of a meaningful vibrating system.

## Conclusions

The present thesis addresses several open issues of theoretical and applicative interest in the field of eigenstructure assignment through active, passive and hybrid control. The main results achieved can be summarized through the following points.

- It has been developed a method for passive control that enables to assign partial mode shapes, in order to attain a better fulfilment of the design specification in the parts of the system whose vibrating behaviour is particularly relevant. The proposed method relies on a homotopy optimization strategy that is proven to boost the achievement of a globally optimal solution of the assignment problem. Such a method has been tested with established benchmarks and in the solution of practical design problems thus proving its effectiveness.
- A novel method for eigenstructure assignment with the addition of auxiliary subsystems has been proposed. For the investigated issue, few methods are available, despite its relevance. The homotopy optimization strategy developed for the assignment of partial eigenvectors have been successfully employed also in this case. The numerical validation of the method proves that the proposed solution is general and reliable and can be applied with different topology of the auxiliary system.
- A method for avoiding spill-over in the assignment of natural frequencies by means of passive control has been also developed. Such an issue can be tackled only for simple systems by the state-of-the-art methods, while the proposed one can handle systems of general topology, even containing distributed parameters, as proved in the thesis.
- The knowledge in the field of the solution of inverse eigenvalue problems has been devoted to another relevant problem in the field of vibrationg system, namely the model updating, which is a field closely related to eigenstructure assignment. The problem has been addressed and a method that enable to perform model updating for flexible link mechanisms is proposed. Such a method has been experimentally applied to a very
challenging updating task, i.e. the tuning of the parameters of a flexible planar robot. by showing that remarkable improvement of the model can be achieved.
- As far as active control is concerned, it has beend developed a method for partial pole placement that enables to deal with spill-over of the non-assigned poles by specifying a regional constraint through Linear Matrix Inequalities. Such a method, in contrast to alternative methods in the literature, employes reliable numerical tools to synthesize the desired control gains. The numerical validation that has been carried out through the challenging test of a flexible link mechanism shows the effectiveness of the proposed approach.
- The capability of assigning the mode shape and the frequency has been significantly enlarged through the last approach proposed in the Thesis, which is a method that exploits passive control to improve the assignment of mode shapes through statefeedback control. Such a hybrid approach has been rarely used in the literature. The proposed method consists in recasting the assignment problem as the minimization of the rank of a matrix, that can be solved exploiting some recent results in the field of numerical optimization. The numerical application has demonstrated the effectiveness of the method and its superiority over the active control and the dynamic structural modification when both used alone and when applied sequentially. In contrast, the proposed method considers both the approaches in a concurrent way.
- Finally, the hybrid method proposed in the thesis has been applied to the experimental control of a cantilever beam by means of a piezoelectric actuator. Despite the fact that the system is highly underactuated, the obtained numerical results proved that the proposed method could achieve a very accurate assignment. Given the impossibility to feed back the whole state vector, the experimental application has imposed to develop a state estimator that enabled to approximate the state of the system even if few measurements are available. The results obtained corroborate all the theoretical expectations and show that the proposed method can be applied to the control of real systems, by taking advantage of state observers suitable for real time control.

The overview of all the results obtained through the different approaches, shows that eigenstructure assignment is an effective tool to improve the dynamic behaviour of vibrating
systems. This statement is corroborated also by the fact that physical feasibility of the computed solution has been taken into account in all the developed methods, as well as their numerical reliability has been carefully addressed. Given the focus on these issues, it is expected that the methods here proposed might be an effective bridge towards the use of eigenstructure assignment in the solution of engineering design problems, such as the mechanical design, the control scheme synthesis, the mechatronics design. With reference to the latter, in particular, the inclusion of the controlled system expected performances in the design of the mechanical part as the hybrid control does, decreases the need of iterating and hence allows to handle effectively more complicated problems.

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